## IN:QUALITIES

G. Hardy, J. E. Littlewood \& G. Pólya

# INEQUALITIES 

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## PREFACE

This book was planned and begun in 1929. Our original intention was that it should be one of the Cambridge Tracts, but it soon became plain that a tract would be much too short for our purpose.

Our objects in writing the book are explained sufficiently in the introductory chapter, but we add a note here about history and bibliography. Historical and bibliographical questions are particularly troublesome in a subject like this, which has applications in every part of mathematics but has never been developed systematically.

It is often really difficult to trace the origin of a familiar inequality. It is quite likely to occur first as an auxiliary proposition, often without explicit statement, in a memoir on geometry or astronomy; it may have been rediscovered, many years later, by half a dozen different authors; and no accessible statement of it may be quite complete. We have almost always found, even with the most famous inequalities, that we have a little new to add.

We have done our best to be accurate and have given all references we can, but we have never undertaken systematic bibliographical research. We follow the common practice, when a particular inequality is habitually associated with a particular mathematician's name; we speak of the inequalities of Schwarz, Hölder, and Jensen, though all these inequalities can be traced further back; and we do not enumerate explicitly all the minor additions which are necessary for absolute completeness.

We have received a great deal of assistance from friends. Messrs G. A. Bliss, L. S. Bosanquet, R. Courant, B. Jessen, V. Levin, R. Rado, I. Schur, L. C. Young, and A. Zygmund have all helped us with criticisms or original contributions. Dr Bosanquet, Dr Jessen, and Prof. Zygmund have read the
proofs, and corrected many inaccuracies. In particular, Chapter III has been very largely rewritten as the result of Dr Jessen's suggestions. We hope that the book may now be reasonably free from error, in spite of the mass of detail which it contains.

Dr Levin composed the bibliography. This contains all the books and memoirs which are referred to in the text, directly or by implication, but does not go beyond them.
G. H. H.
J. E.L.
G. P.

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## CHAPTER I,

## INTRODUCTION

1.1. Finite, infinite, and integral inequalities. It will be convenient to take some particular and typical inequality as a text for the general remarks which occupy this chapter; and we select a remarkable theorem due to Cauchy and usually known as 'Cauchy's inequality'.

Cauchy's inequality (Theorem 7) is

$$
\begin{align*}
& \left(a_{1} b_{1}+a_{2} b_{2}+\ldots+a_{n} b_{n}\right)^{2}  \tag{1.1.1}\\
& \quad \leqq\left(a_{1}{ }^{2}+a_{2}{ }^{2}+\ldots+a_{n}{ }^{2}\right)\left(b_{1}{ }^{2}+b_{2}{ }^{2}+\ldots+b_{n}{ }^{2}\right)
\end{align*}
$$

or

$$
\begin{equation*}
\left(\sum_{1}^{n} a_{\nu} b_{\nu}\right)^{2} \leqq \sum_{1}^{n} a_{\nu}{ }^{2} \sum_{1}^{n} b_{\nu}^{2}, \tag{1.1.2}
\end{equation*}
$$

and is true for all real values of $a_{1}, a_{2}, \ldots, a_{n}, b_{1}, b_{2}, \ldots, b_{n}$. We call $a_{1}, \ldots, b_{1}, \ldots$ the variables of the inequality. Here the number of variables is finite, and the inequality states a relation between certain finite sums. We call such an inequality an elementary or finite inequality.

The most fundamental inequalities are finite, but we shall also be concerned with inequalities which are not finite and involve generalisations of the notion of a sum. The most important of such generalisations are the infinite sums

$$
\begin{equation*}
\sum_{1}^{\infty} a_{\nu}, \sum_{-\infty}^{\infty} a_{\nu} \tag{1.1.3}
\end{equation*}
$$

and the integral

$$
\begin{equation*}
\int_{a}^{b} f(x) d x \tag{1.1.4}
\end{equation*}
$$

(whére $a$ and $b$ may be finite or infinite). The analogues of (1.1.2) corresponding to these generalisations are

$$
\begin{equation*}
\left(\sum_{1}^{\infty} a_{\nu} b_{\nu}\right)^{2} \leqq \sum_{1}^{\infty} a_{\nu} \sum_{1}^{\infty} \sum_{\nu}^{\infty} b^{2} \tag{1.1.5}
\end{equation*}
$$

(or the similar formula in which both limits of summation are infinite), and

$$
\begin{equation*}
\left(\int_{a}^{b} f(x) g(x) d x\right)^{2} \leqq \int_{a}^{b} f^{2}(x) d x \int_{a}^{b} g^{2}(x) d x . \tag{1.1.6}
\end{equation*}
$$

We call (1.1.5) an infinite, and (1.1.6) an integral, inequality.
1.2. Notations. We have often to distinguish between different sets of the variables. Thus in (1.1.2) we distinguish the two sets $a_{1}, a_{2}, \ldots, a_{n}$ and $b_{1}, b_{2}, \ldots, b_{n}$. It is convenient to have a shorter notation for sets of variables, and often, instead of writing 'the set $a_{1}, a_{2}, \ldots, a_{n}$ ' we shall write 'the set (a)' or simply ' the $a$ '.

We shall habitually drop suffixes and limits in summations, when there is no risk of ambiguity. Thus we shall write

$$
\Sigma a
$$

for any of

$$
\sum_{1}^{n} a_{\nu}, \quad \sum_{1}^{\infty} a_{\nu}, \quad \sum_{-\infty}^{\infty} a_{\nu} ;
$$

so that, for example,

$$
\begin{equation*}
(\Sigma a b)^{2} \leqq \Sigma a^{2} \Sigma b^{2} \tag{1.2.1}
\end{equation*}
$$

may mean either of (1.1.2) or (1.1.5), according to the context.
In integral inequalities, the set is replaced by a function; thus in passing from (1.1.2) to (1.1.6), (a) and (b) are replaced by $f$ and $g$. We shall also often omit variables and limits in integrals, writing
$\int f d x$
for (1.1.4): so that (1.1.6), for example, will be written as

$$
\begin{equation*}
\left(\int f g d x\right)^{2} \leqq \int f^{2} d x \int g^{2} d x . \tag{1.2.2}
\end{equation*}
$$

The ranges of the variables, whether in sums or integrals, are prescribed at the beginnings of chapters or sections, or may be inferred unambiguously from the context.
1.3. Positive inequalities. We are interested primarily in 'positive' inequalities ${ }^{\text {a }}$. A finite or infinite inequality is positive if all variables $a, b, \ldots$ involved in it are real and non-negative. An inequality of this type usually carries with it, as a trivial

[^0]corollary, an apparently more general inequality valid for all real, or even complex, $a, b, \ldots$. Thus from (1.1.2) and the inequality
\[

$$
\begin{equation*}
|\Sigma u| \leqq \Sigma|u| \tag{1.3.1}
\end{equation*}
$$

\]

valid for all real or complex $u$, we deduce

$$
\begin{equation*}
|\Sigma a b|^{2} \leqq(\Sigma|a||b|)^{2} \leqq \Sigma|a|^{2} \Sigma|b|^{2} \tag{1.3.2}
\end{equation*}
$$

where the $a$ and $b$ are arbitrary complex numbers. We shall usually be content to state our theorems in the fundamental 'positive' form and to leave the derived results to the reader. Occasionally, however, when the inequality in question is very important, we state it explicitly in its most general form.

Similar remarks apply to integral inequalities. The independent variable $x$ will be real, but will (like the variable of summation $\nu$ ) take either positive or negative values; while the functions $f(x)$, $g(x), \ldots$ will generally assume non-negative values only. To such on inequality as (1.1.6), true for non-negative $f, g$, corresponds the more general inequality
(1.3.3) $\quad\left|\int f g d x\right|^{2} \leqq \int|f|^{2} d x \int|g|^{2} d x$, valid for arbitrary complex functions $f, g$ of the real variable $x$.

Numbers $k, l, r, s, \ldots$ occurring as indices in our theorems are seail but in general capable of either sign.
1.4. Homogeneous inequalities. The two sides of (1.1.2) are homogeneous functions of degree 2 of the $a$ and also of the $b$; and generally both sides of our inequalities will be homogeneous functions, of the same degree, of certain sets of variables. Since homogeneous functions of positive degree vanish when all their arguments vanish, both sides, if of positive degree, will vanish, and so be equal, when the sets concerned consist entirely of 0 's. Thus (1.1.2) reduces to an equality if all the $a$, or all the $b$, are 0 .

- A set consisting entirely of 0's is called a nul set, or the nul set, if the context is unambiguous. In general the ' $\leqq$ ' or ' $\geqq$ ' of our -theorems will reduce to ' $=$ ' when one or all of the sets involved is nul. Sometimes this will be the only case of equality. More usually there will be other cases; thus plainly ' $=$ ' occurs in (1.1.2) if every $a$ is equal to the corresponding $b$. We shall be careful, wherever it is possible, to pick out explicitly such cases of equality.

The homogeneity of an inequality in certain sets of variables often enables us to simplify our proofs by imposing an additional restriction (a normalisation) on them. Thus the means $\mathfrak{M}_{r}(a)$ of $\S 2.2$ are homogeneous, of degree 0 , in the weights $p$, and we may always suppose, if we please, that $\Sigma p=1$. Again, if we wish to prove that

$$
\begin{equation*}
\left.\left(a_{1}^{s}+a_{2}^{s}+\ldots+a_{n}\right)^{s}\right)^{1 / s} \leqq\left(a_{1}^{r}+a_{2}^{r}+\ldots+a_{n}^{r}\right)^{1 / r} \tag{1.4.1}
\end{equation*}
$$

when $0<r<s$ (Theorem 19), we may suppose (since both sides are homogeneous in the $a$ of degree 1) that $\Sigma a^{r}=1$. We have then

$$
a_{\nu}{ }^{r} \leqq 1, \quad a_{\nu}^{s}=\left(a_{\nu}^{r}\right)^{s / r} \leqq a_{\nu}{ }^{r},
$$

and so $\Sigma a^{s} \leqq \Sigma a^{r}=1$. Without this preliminary normalisation, our proof would run

$$
\frac{\left(\Sigma a^{s}\right)^{1 / s}}{\left(\Sigma a^{r}\right)^{1 / r}}=\left\{\Sigma \frac{a^{s}}{\left(\Sigma a^{r}\right)^{s / r}}\right)^{1 / s}=\left\{\Sigma\left(\frac{a^{r}}{\Sigma a^{r}}\right)^{s / r}\right\}^{1 / s} \leqq\left(\Sigma \frac{a^{r}}{\Sigma a^{r}}\right)^{1 / s}=1 .
$$

There is another sense of 'homogeneity' which is sometimes important. Let us compare (1.4.1) above, which may be written as

$$
\begin{equation*}
\left(\mathrm{\Sigma} a^{s}\right)^{1 / s} \leqq\left(\mathrm{\Sigma} a^{r}\right)^{1 / r}, \tag{1.4.2}
\end{equation*}
$$

with (1.1.2). Both inequalities are homogeneous in the variables, but (1.1.2) has a further homogeneity which (1.4.2) has not. It is, as we may say, 'homogeneous in $\Sigma$ '; $\Sigma$, if treated as a number, would occur to the same power on the two sides of the inequality.

The result of this homogeneity in $\Sigma$ is that (1.1.2) remains true if every sum which occurs is replaced by the corresponding mean, i.e. if written in the form

$$
\cdot\left(\frac{1}{n} \Sigma a b\right)^{2} \leqq\left(\frac{1}{n} \Sigma a^{2}\right)\left(\frac{1}{n} \Sigma b^{2}\right) .
$$

The importance of this kind of homogeneity will appear very clearly in § 2.10 and § 6.4. Roughly, an inequality which possesses it has an integral analogue, while one which does not, like (1.4.2), has none.
1.5. The axiomatic basis of algebraic inequalities ${ }^{\text {a }}$. Our subject is difficult to define precisely, but belongs partly to 'algebra' and partly to 'analysis'. Algebra or analysis, like geometry, may be treated axiomatically. Instead of saying, as

[^1]for example in Dedekind's theory of real numbers, that we are concerned with such or such definite objects, we may say, as in projective geometry, that we are concerned with any system of objects which possesses certain properties specified in a set of axioms. We do not propose to consider the 'axiomatics' of different parts of the subject in detail, but it may be worth while to insert a few remarks concerning the axiomatic basis of those theorems which, like (1.1.2) and most of the theorems of Ch. II, belong properly to algebra.

We may take as the axioms of an algebra only the ordinary laws of addition and multiplication. All our theorems will then be true in many different fields, in real algebra, complex algebra, or the arithmetic of residues to any modulus. Or we may add axioms concerning the solubility of linear equations, axioms which secure the existence and uniqueness of difference and quotient. Our theorems will then be true in real or complex algebra or in arithmetic to a prime modulus.

In our present subject we are concerned with relations of inequality, a notion peculiar to real algebra. We can secure an axiomatic basis for theorems of inequality by taking, in addition to the 'indefinables' and axioms already referred to, one new indefinable and two new axioms. We take as indefinable the idea of a positive number, and as axioms the two propositions:
I. Either $a$ is 0 or $a$ is positive or $-a$ is positive, and these possibilities are exclusive.
II. The sum and product of two positive numbers are positive.

We say that $a$ is negative if $-a$ is positive, and that $a$ is greater (less) than $b$ if $a-b$ is positive (negative). Any inequality of a purely algebraic type, such as (1.1.2), may be made to rest on this foundation.
1.6. Comparable functions. We may say that the functions

$$
f(a)=f\left(a_{1}, a_{2}, \ldots, a_{n}\right), \quad g(a)=g\left(a_{1}, a_{2}, \ldots, a_{n}\right)
$$

are comparable if there is an inequality between them valid for all non-negative real $a$, that is to say if either $f \leqq g$ for all such $a$ or
$f \geqq g$ for all such $a$. Two given functions are not usually comparable. Thus two positive homogeneous polynomials of different degrees are certainly not comparable ; if $0 \leqq f \leqq g$ for all nonnegative $a$, and both sides are homogeneous, then $f$ and $g$ are certainly of the same degree.

The definition may naturally be extended to functions $f(a, b, \ldots)$ of several sets of variables.

We shall be occupied throughout this volume with problems concerning the comparability of functions. Thus the arithmetic and geometric means of the $a$ are comparable: $\mathscr{E}(a) \leqq \mathfrak{A}(a)$ (Theorem 9). The functions $\mathfrak{5}(a+b)$ and $\mathfrak{5}(a)+(\mathscr{5}(b)$ are comparable (Theorem 10). The functions $\mathfrak{A}(a b)$ and $\mathfrak{H}(a) \mathfrak{X}(b)$ are not comparable; their relative magnitude depends upon the relations of magnitude of the $a$ and $b$ (Theorem 43). The functions

$$
\psi^{-1}(\Sigma p \psi(a)), \quad \chi^{-1}(\Sigma p \chi(a))
$$

are comparable if and only if $\chi \psi^{-1}$ is convex or concave (Theorem 85).

An important general theorem concerning the comparability of two functions of the form

$$
\Sigma a_{1}{ }^{\alpha_{1}} a_{2}{ }^{\alpha_{2}} \ldots a_{n}^{\alpha_{n}},
$$

due to Muirhead, will be found in §2.18.
1.7. Selection of proofs. The methods of proof which we use in different parts of the book will depend on very different sets of ideas, and we shall often, particularly in Ch. II, give a number of alternative proofs of the same theorem. It may be useful to call attention here to certain broad distinctions between the methods which we employ.

In the first place, many of the proofs of Ch. II are 'strictly elementary', since they depend solely on the ideas and processes of finite algebra. We have made it a principle to give at any rate one such proof of any really important theorem whose character permits it.

Next we have, even in Ch. II, many proofs which are not elementary in this sense because they involve considerations of

[^2]limits and continuity. We have also, particularly in Ch. IV, proofs which depend upon the standard properties of differential coefficients, as for example upon Rolle's Theorem. All these proofs belong to the elements of the theory of functions of one real variable.

Later, when dealing with integrals in Ch. VI, we naturally make use of the theory of measure and of the integral of Lebesgue. This we take for granted, but we give a summary in §§ 6.1-6.3 of the parts of the theory which we require.

Occasionally we appeal to the more remote parts of the theory of functions of real variables; but we do this only in alternative proofs or in the proofs of theorems of considerable intrinsic difficulty. Thus in Ch.IV (§4.6) we use the theory of the maxima and minima of functions of several variables; in Ch. VII we use the methods of the Calculus of Variations; and in Ch. IX we use the theory of double and repeated integration. We make no use of complex function theory, although, in the last chapters, we refer to it occasionally for purposes of illustration. The sections in which we do this do not belong properly to the main body of the book.

We add a few further remarks of a more detailed character.
(i) Cauchy's inequality (1.1.2) is a proposition of finite algebra, as defined in §1.5. It is a recognised principle that the proof of such a theorem should involve only the methods of the theory to which it belongs.
(ii) We shall be continually meeting theorems, such as Hölder's inequality

$$
\begin{equation*}
\Sigma a b \leqq\left(\Sigma a^{k}\right)^{1 / k}\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}} \tag{1.7.1}
\end{equation*}
$$

(Theorem 13), whose status depends uponthe value of a parameter $k$. If $k$ is rational, the theorem is algebraical, and our remarks under (i) apply. If $k$ is irrational, $a^{k}$ is not an algebraical function, and it is obvious that there can be no strictly algebraical proof.

It is however reasonable to demand, when we are concerned with an inequality so fundamental as Hölder's, that our step outside algebra shall be the absolute minimum which the nature
of the problem necessitates. It is plain that this step will depend upon our definition of $a^{k}$. We may define $a^{k}$ as $\exp (k \log a)$, and in this case it is obviously legitimate and necessary to use the theory of the exponential and logarithmic functions. If, as is more usual, we define $a^{k}$ as the limit of $a^{k_{n}}$, where $k_{n}$ is an appropriate rational approximation to $k$, then this limiting process should be the only one to which we appeal.
(iii) Suppose that, adopting the last point of view, we have proved Hölder's inequality, for rational $k$, in the form (1.7.1.). We can infer its truth for irrational $k$ by a passage to the limit.

Such a proof, however, is not usually sufficient for our purpose. We always wish to prove a theorem of a more precise type than (1.7.1), in which (as in Theorem 13) we establish strict inequality except in certain specified special cases. When we pass to the limit, ' <' becomes ' $\leqq$ ', we lose touch with the cases of equality (though these are in fact the same as in the rational case), and our proof is incomplete. It is therefore necessary to arrange our proofs in such a manner as to avoid such passages to the limit wherever it is possible. The same point arises whenever we wish to pass from a finite inequality to the corresponding infinite or integral inequality. It recurs at intervals throughout the volume and has often determined our choice of a particular line of proof.
(iv) The general principles which have governed our choice of methods are as follows. When a theorem is simple and fundamental, like Theorems 7, 9 , or 11, we prove it by several different methods, and are careful that one of our methods at any rate shall conform to the canons laid down under (i) and (ii). When the theorem is subsidiary or difficult, or when a proof satisfying these conditions would be troublesome orlong, we use whatever method seems to us simplest or most instructive.
1.8. Selection of subjects. The principles which have guided us in our selection of subjects may be summarised as follows.
(i) The first part of the book (Chs. II-VI ${ }^{\text {a }}$ ) contains a systematic treatment of a definite subject. Our object has been to

[^3]discuss thoroughly (with their analogues and extensions) the simple inequalities which are 'in daily use' in analysis. Of these three are fundamental, viz.
(1) the theorem of the arithmetic and geometric means (Theorem 9),
(2) Hölder's inequality (Theorem 11),
(3) Minkowski's inequality (Theorem 24);
and these three theorems dominate the first six chapters. We prove them in a variety of ways, in the finite case in Ch. II, in the infinite case in Ch. V, and in the integral case in Ch. VI; while Ch. III (which contains a general account of the theory of convex functions) is mainly occupied with their generalisations. In these chapters, of which the most important are II, III, and VI, we have aimed at a comprehensive and in some ways exhaustive treatment.
(ii) The rest of the book (Chs. VII-X) is written in a different spirit and must be judged by different standards. These chapters contain a series of essays on subjects suggested by the more systematic investigations which precede. In them there is very little attempt at system or completeness. They are intended as an introduction to certain fields of modern research, and we have allowed our personal interests to dominate our choice of topics.

In spite of this (or because of it) the chapters have a certain unity. There is much modern work, in real or complex functiontheory, in the theory of Fourier series, or in the general theory of orthogonal developments, in which the 'Lebesgue classes $L^{k}$ ' occupy the central position. This work demands a considerable mastery of the technique of inequalities; Hölder's and Minkowski's inequalities, and other more modern and more sophisticated inequalities of the same general character, are required at every turn. Our object has been to write such an introduction to this field of analysis as may be made to hang naturally on the subject matter of the early chapters.
(iii) We are interested primarily in certain parts of real analysis, and not in arithmetic or in algebra for its own sake. The line
between algebra and analysis is often difficult to draw, especially in the theory of quadratic or bilinear forms, and we have often doubted what to include or reject. We have however excluded all developments whose main interest seemed to us to be algebraical.

We have also excluded function-theory proper, real or complex. In the later chapters, however, we have sometimes tried to show the significance of our theorems by sketching the lines of some of their function-theoretic applications.

Thus (to give definite examples) our programme excludes
(1) inequalities of a definitely arithmetical character, such as those of the theory of primes, or those which give bounds for forms with integral variables;
(2) inequalities which belong properly to the algebraical theory of quadratic forms;
(3) inequalities, such as 'Bessel's inequality', which belong to the theory of orthogonal series;
(4) inequalities, such as 'Hadamard's three circle theorem', which belong to function-theory proper:
and there is no systematic discussion of geometrical inequalities, though we use them frequently for purposes of illustration.

It may be useful to end this introduction by a few words of advice to readers who are anxious to avoid unnecessary immersion in detail. The subject, attractive as it is, demands, for the writer at any rate, a great deal of attention to details of a rather tiresome kind. These details arise particularly in the exclusion of exceptional cases, the complete specification of cases of equality, and the conventional treatment of zero and infinite values. Such a reader as we have in mind may be content, in general, to simplify his task as follows. (1) He may ignore the distinction between non-negative and positive, so that the numbers and functions with which he is concerned are all positive in the narrow sense. (2) He may ignore our conventions concerning 'infinite values'. (3) He may assume that the parameter $k$ or $r$ of inequalities such as Hölder's and Minkowski's is greater than 1.
(4) He may take it for granted that 'what goes for sums goes, with the obvious modifications, for integrals' (or vice versa). He should then be able to master what is essential without undue trouble.

This advice for 'easy reading' must not be taken too literally. It is essential to understand the kind of exceptional cases which occur, and the general principles which govern the discrimination of cases of equality. It is not a mere academic exercise to pick out the cases of equality in such an inequality as Hölder's; a knowledge of these cases provides (as is shown very clearly in $\S \S 8.13-8.16$ ) a powerful weapon for the discovery of deep and important theorems. Every reader should make it his business to explore this inequality at any rate to the end.

## CHAPTER II

## ELEMENTARY MEAN VALUES

2.1. Ordinary means. In what follows we are concerned with sets of $n$ non-negative numbers $a$ (or $b, c, \ldots$ ), say

$$
\begin{equation*}
a_{1}, a_{2}, \ldots, a_{\nu}, \ldots, a_{n} \quad\left(a_{\nu} \geqq 0\right), \tag{2.1.1}
\end{equation*}
$$

and a real parameter $r$, which we suppose for the present not to be zero.

We denote the ordered series (2.1.1) by (a). When we say that ' $(a)$ is proportional to (b)' we mean that there are two numbers $\lambda$ and $\mu$, not both zero, such that

$$
\begin{equation*}
\lambda a_{\nu}=\mu b_{\nu} \quad(\nu=1,2, \ldots, n) . \tag{2.1.2}
\end{equation*}
$$

It will be observed that the nul set, the set ( $a$ ) in which every $a$ is zero, is proportional to any (b). Proportionality, as we have defined it, is a symmetrical relation between sets but not a transitive one; it becomes transitive if we exclude the nul set from consideration.
If (a) and (b) are proportional, and neither of them is nul, then $b_{v}=0$ whenever $a_{\nu}=0$, and $a_{\nu} / b_{\nu}$ is independent of $\nu$ for the remaining values of $\nu$.

We write

$$
\begin{equation*}
\mathfrak{M}_{r}=\mathfrak{M}_{r}(a)=\left(\frac{1}{n} \Sigma a^{r}\right)^{1 / r}=\left(\frac{1}{n} \sum_{v=1}^{n} a_{v}^{r}\right)^{1 / r}, \tag{2.1.3}
\end{equation*}
$$

except when (i) $r=0$ or (ii) $r<0$ and one or more of the $a$ are zero. In the exceptional case (ii), when (2.1.3) has no meaning, we define $\mathfrak{M}_{r}$ as zero, so that

$$
\begin{equation*}
\mathfrak{M}_{r}=0 \quad(r<0, \text { some a zero })^{\mathrm{a}} . \tag{2.1.4}
\end{equation*}
$$

Here and elsewhere we shall omit the suffixes and limits of summation when it can be done without ambiguity.

In particular we write

$$
\begin{align*}
& \mathfrak{A}=\mathfrak{A}(a)=\mathfrak{M}_{1}(a),  \tag{2.1.5}\\
& \mathfrak{G}=\mathfrak{G}(a)=\mathfrak{M}_{-1}(a) . \tag{2.1.6}
\end{align*}
$$

Finally, we write

$$
\left.\mathfrak{G}=\mathfrak{G}(a)=\sqrt[n]{\left(a_{1}\right.} a_{2} \ldots a_{n}\right)=\sqrt[n]{(\Pi a) .}
$$

[^4]Thus $\mathfrak{A}(a), \mathfrak{G}(a)$, $\mathscr{S}(a)$ are the ordinary arithmetic, harmonic, and geometric means.

We have excluded the case $r=0$, but we shall find later (§ 2.3) that we can interpret $\mathfrak{M}_{0}$ conventionally as $\mathfrak{F}$. We are not generally concerned with negative $a$, but it is sometimes convenient to use $\mathfrak{A}(a)$ without any restriction of sign. The definition is unchanged.
2.2. Weighted means. We shall however usually work with a more general system of mean values. We suppose that

$$
\begin{equation*}
p_{\nu}>0 \quad(\nu=1,2, \ldots, n) \tag{2.2.1}
\end{equation*}
$$

and write

$$
\begin{gather*}
\mathfrak{M}_{r}=\mathfrak{M}_{r}(a)=\mathfrak{M}_{r}(a, p)=\left(\frac{\Sigma p a^{r}}{\Sigma p}\right)^{1 / r},  \tag{2.2.2}\\
\mathfrak{M}_{r}=0 \quad(r<0, \text { some a zero }), \tag{2.2.3}
\end{gather*}
$$

$$
\begin{equation*}
\mathfrak{( S}=\mathfrak{G}(a)=\mathfrak{G}(a, p)=\left(\Pi a^{p}\right)^{1 / \Sigma p} . \tag{2.2.4}
\end{equation*}
$$

The equations (2.1.5) and (2.1.6) stand as before, with the addition of the symbols $\mathfrak{A}(a, p), \mathfrak{G}(a, p)$. The last remark of $\S 2.1$ applies also to the generalised $\mathfrak{A}$. The weighted means reduce to the ordinary means when $p_{\nu}=1$ for every $\nu$.

The means being homogeneous and of degree 0 in the $p$ 's, we may suppose, if we please, that $\Sigma p=1$. In this case we shall replace $p$ by $q$; thus

$$
\begin{array}{cc}
\mathfrak{M}_{r}(a)=\mathfrak{M}_{r}(a, q)=\left(\Sigma q a^{r}\right)^{1 / r} & (\Sigma q=1) \\
\mathfrak{G S}(a)=\mathfrak{S}(a, q)=\Pi a^{q} & (\Sigma q=1) \tag{2.2.6}
\end{array}
$$

We shall not usually refer to the weights explicitly in our formulae, but it is always to be understood that mean values which are compared with one another are formed with the same weights.

Ordinary means are special cases of weighted means. On the other hand, weighted means with commensurable weights are special cases of ordinary means (with a different system of $a$ ); for we may suppose, on account of homogeneity, that the weights are integral, and we can derive means with integral weights from ordinary means by replacing every number by an appropriate set of equal numbers. Means with incommensurable weights may be regarded as limiting cases of ordinary means.

The following obvious formulae will be used repeatedly:

$$
\begin{align*}
\mathfrak{M}_{r}(a) & =\left\{\mathfrak{H}\left(a^{r}\right)\right\}^{1 / r},  \tag{2.2.7}\\
\mathfrak{G}(a) & =e^{\mathfrak{H}(\log a)}
\end{align*}
$$

$$
\begin{align*}
& \mathfrak{M}_{-r}(a)=\frac{1}{\mathfrak{M}_{r}(1 / a)}  \tag{2.2.9}\\
& \mathfrak{M}_{r s}(a)=\left\{\mathfrak{M}_{s}\left(a^{r}\right)\right\}^{1 / r} . \tag{2.2.10}
\end{align*}
$$

We suppose that $a>0$ in (2.2.8), and in the other formulae if a suffix is negative; the formulae may be extended to cover the missing cases by appropriate conventions. Also

$$
\begin{gather*}
\mathfrak{H}(a+b)=\mathfrak{A}(a)+\mathfrak{H}(b),  \tag{2.2.11}\\
\mathfrak{G}(a b)=\mathfrak{S}(a) \mathfrak{G}(b),  \tag{2.2.12}\\
\mathfrak{M}_{r}(b)=k \mathfrak{M}_{r}(a) \quad \text { if }(b)=k(a) \tag{2.2.13}
\end{gather*}
$$

(i.e. if $b_{\nu}=k a_{\nu}$, where $k$ is independent of $\nu$ ),

$$
\begin{gather*}
\mathfrak{S}(b)=k \mathscr{( S}(a) \quad \text { if }(b)=k(a),  \tag{2.2.14}\\
\mathfrak{M}_{r}(a) \leqq \mathfrak{M}_{r}(b) \quad \text { if } a_{\nu} \leqq b_{\nu} \text { for all } \nu . \tag{2.2.15}
\end{gather*}
$$

2.3. Limiting cases of $\mathfrak{M}_{r}(a)$. We denote by

$$
\operatorname{Min} a, \quad \operatorname{Max} a
$$

the smallest and largest value of an $a$.

1. $\operatorname{Min} a<\mathfrak{M}_{r}(a)<\operatorname{Max} a$, unless either all the a are equal, or else $r<0$ and an a is zero.

It is to be understood here, and in the enunciations of all later theorems, that, when we assert that inequalities hold unless some particular condition is satisfied, we imply that at least one of the inequalities degenerates into an equality in the case excluded. Here, for example, $\operatorname{Min} a=\mathfrak{M d}_{r}(a)=\operatorname{Max} a$ if all $a$ are equal, and $\operatorname{Min} a=\mathbb{M}_{r}(a) \leqq \operatorname{Max} a$ in the other exceptional case.

We form our means with $q$. Since

$$
\Sigma q(a-\mathfrak{Y})=0
$$

every $a$ is equal to $\mathfrak{A}$, or else $a-\mathfrak{A}$ is positive for at least one $a$ and negative for another. This proves the theorem for $r=1$.

In the general case we may suppose that either $a>0$ or else $r>0$, the cases excluded being trivial. It then follows that

$$
\left\{\mathfrak{M} \mathcal{l}_{r}(a)\right\}^{r}=\mathfrak{A}\left(a^{r}\right)
$$

lies between (Mina) and (Maxa)r, which proves the theorem generally.
2. $\operatorname{Min} a<\mathscr{G}(a)<\operatorname{Max} a$, unless all the $a$ are equal or an $a$ is zero. In the second exceptional case $\mathfrak{F}=0$. If $\mathfrak{F}>0$ then

$$
\Pi\left(\frac{a}{\sqrt{6}}\right)^{q}=1,
$$

so that every $a$ is $\mathfrak{F S}$ or at least one is greater and one less than $\mathscr{F}$.
3. $\lim _{r \rightarrow 0} \mathfrak{M}_{r}(a)=\mathscr{F}(a)$.

If every $a$ is positive

$$
\begin{aligned}
\mathfrak{M}_{r}(a) & =\exp \left(\frac{1}{r} \log \Sigma q a^{r}\right) \\
& =\exp \left\{\frac{1}{r} \log \left(1+r \Sigma q \log a+O\left(r^{2}\right)\right)\right\} \\
& \rightarrow \exp (\Sigma q \log a)=\Pi a^{q}=\mathscr{S}(a),
\end{aligned}
$$

when $r \rightarrow 0$.
If there are some zero $a, b$ denotes a positive $a$, and $s$ is a $q$ corresponding to a $b$, then

$$
\mathfrak{M}_{r}(a, q)=\left(\Sigma q a^{r}\right)^{1 / r}=\left(\Sigma s b^{r}\right)^{1 / r}=(\Sigma s)^{1 / r} \mathfrak{M}_{r}(b, s) \rightarrow 0
$$

when $r \rightarrow+0$, since $\mathfrak{M}_{r}(b, s) \rightarrow \mathfrak{F}(b, s)$ and $\Sigma s<1$. When $r<0, \mathfrak{M}_{r}$ and (5) are both zero, so that the result holds also when $r \rightarrow-0$.

Our proof depends on the theory of the exponential and logarithmic functions. We show in $\S 2.16$ how a more elementary proof may be found if desired.
4. $\lim _{r \rightarrow \infty} \mathfrak{M}_{r}(a)=\operatorname{Max} a, \lim _{r \rightarrow-\infty} \mathfrak{M}_{r}(a)=\operatorname{Min} a$.

If $a_{k}$ is the largest $a$, or one of the largest, and $r>0$, we have

$$
q_{k}{ }^{1 / r} a_{k} \leqq \mathfrak{M}_{r}(a) \leqq a_{k} ;
$$

from which the first equation follows at once. The second is trivial if any $a$ is zero and follows from (2.2.9) otherwise.

We now agree to write
(2.3.1) $\quad \mathfrak{M}_{0}(a)=\mathfrak{G}(a), \quad \mathfrak{M}_{\infty}(a)=\operatorname{Max} a, \quad \mathfrak{R}_{-\infty}(a)=\operatorname{Min} a$.

With these conventions, we have
5. $\mathfrak{M}_{-\infty}(a)<\mathfrak{M}_{r}(a)<\mathfrak{M}_{\infty}(a)$ for all finite $r$, unless the a are all equal, or $r \leqq 0$ and an a is zero.
2.4. Cauchy's inequality. It is convenient to prove the next theorem here although it will be superseded later by a more complete theorem (Theorem 16).
6. $\mathfrak{M}_{r}(a)<\mathfrak{M}_{2 r}(a)(r>0)$, unless all the $a$ are equal.

The inequality is

$$
\left(\Sigma p a^{r}\right)^{2}<\Sigma p \Sigma p a^{2 r}
$$

and is a special case of the very important theorem which follows.
7. $(\Sigma a b)^{2}<\Sigma a^{2} \Sigma b^{2}$, unless (a) and (b) are proportionala.

For

$$
\Sigma a^{2} \Sigma b^{2}-(\Sigma a b)^{2}=\frac{1}{2} \sum_{\mu, \nu}\left(a_{\mu} b_{\nu}-a_{\nu} b_{\mu}\right)^{2} .
$$

An alternative proof is as follows. The quadratic form

$$
\Sigma(x a+y b)^{2}=x^{2} \Sigma a^{2}+2 x y \Sigma a b+y^{2} \Sigma b^{2}
$$

is positive for all $x, y$, and therefore has a negative discriminant, unless $x a_{\nu}+y b_{\nu}=0$ for some $x, y$, not both zero, and all $\nu$.

To deduce Theorem 6, take $\sqrt{ } p$ and $a^{r} \sqrt{ } p$ in place of $a$ and $b$.
Theorem 7 may be generalised as follows:
8.

$$
\left.\begin{array}{cccc|}
\Sigma a^{2} & \Sigma a b & \ldots & \Sigma a l \\
\ldots & \ldots & \ldots & \ldots \\
\Sigma l a & \Sigma l b & \ldots & \Sigma l^{2}
\end{array} \right\rvert\,>0,
$$

unless the sets $(a),(b), \ldots,(l)$ are linearly dependent, i.e. unless there are numbers $x, y, \ldots, w$, not all zero, such that $x a_{\nu}+y b_{\nu}+\ldots+w l_{\nu}=0$ for every $\nu$.

Either proof of Theorem 7 may be extended to prove Theorem 8: we may either express the determinant as a sum of squares of determinants, or we may consider the non-negative quadratic form

$$
\Sigma(x a+y b+\ldots+w l)^{2}
$$

in $x, y, \ldots, w$. We do not go into details because any systematic discussion of inequalities connected with determinants and quadratic forms would carry us beyond the limits which we have imposed on the book.
2.5. The theorem of the arithmetic and geometric means. We come now to the most famous theorem of the subject.

[^5]9. $\mathfrak{G}(a)<\mathfrak{A}(a)$, unless all the a are equal.

The inequality to be proved may be written in either of the forms

$$
\begin{equation*}
a_{1}{ }^{p_{1}} a_{2}^{p_{2}} \ldots a_{n}^{p_{n}}<\left(\frac{p_{1} a_{1}+\ldots+p_{n} a_{n}}{p_{1}+\ldots+p_{n}}\right)^{p_{1}+\ldots+p_{n}} \tag{2.5.1}
\end{equation*}
$$

$$
\begin{equation*}
a_{1}{ }_{1}^{q_{1}} a_{2}^{q_{2}} \ldots a_{n}{ }^{q_{n}}<\Sigma q a \tag{2.5.2}
\end{equation*}
$$

(where as usual $\Sigma q=1$ ).
This theorem is so fundamental that we propose to give a pumber of proofs, of varying degrees of simplicity and generality. Of the two which we give in this section, the first is entirely elementary. The second depends on Theorem $\mathbf{3}$ and so, at present, on the theory of the exponential and logarithmic functions. We shall show later ( $\S 2.16$ ) how this proof also may be made to conform more strictly to the canons of § 1.7.
(i) ${ }^{\text {a }}$ We have ${ }^{\text {b }}$

$$
a_{1} a_{2}=\left(\frac{a_{1}+a_{2}}{2}\right)^{2}-\left(\frac{a_{1}-a_{2}}{2}\right)^{2}<\left(\frac{a_{1}+a_{2}}{2}\right)^{2},
$$

unless $a_{1}=a_{2}$, and so

$$
a_{1} a_{2} a_{3} a_{4} \leqq\left(\frac{a_{1}+a_{2}}{2}\right)^{2}\left(\frac{a_{3}+a_{4}}{2}\right)^{2} \leqq\left(\frac{a_{1}+a_{2}+a_{3}+a_{4}}{4}\right)^{4},
$$

with inequality in one place or the other unless $a_{1}=a_{2}=a_{3}=a_{4}$. Repeating the argument $m$ times, we find

$$
\begin{equation*}
a_{1} a_{2} \ldots a_{2^{m}}<\left(\frac{a_{1}+a_{2}+\ldots+a_{2^{m}}}{2^{m}}\right)^{2^{m}}, \tag{2.5.3}
\end{equation*}
$$

unless all the $a$ are equal. This is (2.5.1) with unit weights and $n$ a power of 2 .

Suppose now that $n$ is any number less than $2^{m}$. Taking

$$
\begin{gathered}
b_{1}=a_{1}, \quad b_{2}=a_{2}, \ldots, b_{n}=a_{n}, \\
b_{n+1}=b_{n+2}=\ldots=b_{2^{m}}=\frac{a_{1}+a_{2}+\ldots+a_{n}}{n}=\mathfrak{A},
\end{gathered}
$$

and applying (2.5.3) to the $b$, we find

$$
a_{1} a_{2} \ldots a_{n} \mathfrak{H}^{2^{m}-n}<\left(\frac{b_{1}+b_{2}+\ldots+b_{2^{m}}}{2^{m}}\right)^{2^{m}}=\left(\frac{n \mathfrak{A}+\left(2^{m}-n\right) \mathfrak{A}}{2^{m}}\right)^{2^{m}}=\mathfrak{A}^{2 m},
$$

or

$$
a_{1} a_{2} \ldots a_{n}<\mathfrak{A}^{n}
$$

$$
\text { a Cauchy }(1,375) . \quad \text { b Euclid }(1: \text { II } 5, \text { v 25) }
$$

unless all the $b$, and so all the $a$, are equal. This is (2.5.1) with unit weights. We deduce (2.5.1), with any commensurable weights, by the process explained in §2.2.

When the weights are incommensurable, we can replace them by a set of commensurable approximations, prove (2.5.1) with the approximating weights, and proceed to the limit. In this process ' $<$ ' is changed into ' $\leqq$ ', so that we do not at first obtain a complete proof of the theorem. We may complete the proof as follows. Write

$$
q_{\nu}=q_{\nu}^{\prime}+q_{\nu}^{\prime \prime} \quad(\nu=1,2, \ldots, n)
$$

where $q_{\nu}^{\prime}>0, q_{\nu}^{\prime \prime}>0$, and $q_{\nu}^{\prime}$ is rational. Then

$$
r^{\prime}=\Sigma q_{\nu}^{\prime}, \quad r^{\prime \prime}=\Sigma q_{\nu}^{\prime \prime}
$$

are rational and $r^{\prime}+r^{\prime \prime}=1$. We have already proved (2.5.1) with ' $<$ ' for rational $p$, and with ' $\leqq$ ' in any case. Hence

$$
\begin{aligned}
\Pi a^{q^{\prime}} & <\left(\frac{\Sigma q^{\prime} a}{\Sigma q^{\prime}}\right)^{\Sigma q^{\prime}}, \quad \Pi a^{q^{\prime \prime}} \leqq\left(\frac{\Sigma q^{\prime \prime}}{\Sigma q^{\prime \prime}}\right)^{\Sigma q^{\prime \prime}}, \\
\Pi a^{q} & =\Pi a^{q^{\prime}} \Pi a^{q^{\prime \prime}}<\left(\frac{1}{r^{\prime}} \Sigma q^{\prime} a\right)^{r^{\prime}}\left(\frac{1}{r^{\prime \prime}} \Sigma q^{\prime \prime} a\right)^{r^{\prime \prime}} \\
& \leqq \Sigma q^{\prime} a+\Sigma q^{\prime \prime} a=\Sigma q a .
\end{aligned}
$$

Another way of completing the proof was shown us by R. E. A. C. Paley. This depends on Theorem 6. From this theorem, the formula (2.2.10), and what has been proved before, it follows that

$$
\mathfrak{A}(a)=\mathfrak{M}_{1}(a)>\mathfrak{M}_{\frac{1}{2}}(a)=\mathfrak{M}_{1}^{2}\left(a^{\frac{1}{2}}\right) \geqq \mathfrak{G}^{2}\left(a^{\frac{1}{2}}\right)=\mathfrak{G}(a) .
$$

(ii) ${ }^{\text {a }}$ By Theorems 6 and 3, we have

$$
\mathfrak{A}(a)=\mathfrak{M}_{1}(a)>\mathfrak{M}_{\frac{1}{2}}(a)>\mathfrak{M}_{\frac{1}{4}}(a)>\ldots>\lim _{m \rightarrow \infty} \mathfrak{M}_{2^{-m}}(a)=\mathfrak{J}(a) .
$$

This proof is very concise but not quite so elementary as the first. It may be observed that we require Theorem 3 only in the case in which the $r$ of Theorem 3 tends to zero through the special sequence of values $2^{-m}$.
2.6. Other proofs of the theorem of the means. We shall return to Theorem 9 in §§ 2.14-15 and again in §2.21. We add here a few remarks about alternative proofs of the ordinary form of the theorem with unit weights.
(i) ${ }^{\text {a }}$ If the $a$ are not all equal, let

$$
a_{1}=\operatorname{Min} a<\operatorname{Max} a=a_{2} .
$$

If we replace each of $a_{1}$ and $a_{2}$ by $\frac{1}{2}\left(a_{1}+a_{2}\right), \mathfrak{A}(a)$ is unaltered, but

$$
\left(\frac{a_{1}+a_{2}}{2}\right)^{2}>a_{1} a_{2}
$$

so that $\mathfrak{F}(a)$ is increased.
Suppose now that we vary the $a$ in such a manner that $\mathfrak{A}$ is constant, and that we assume the existence of a set ( $a^{*}$ ) for which (8) attains a maximum value. Then the $a^{*}$ must be equal, since if not we can replace them as above by another system for which $\mathscr{O}_{6}$ is greater. It follows that the maximum of $\mathfrak{f}$ is $\mathfrak{A}$, and that this maximum is attained only for equal $a$.

To prove the existence of ( $a^{*}$ ), let

$$
\phi\left(a_{1}, a_{2}, \ldots, a_{n-1}\right)=a_{1} a_{2} \ldots a_{n-1}\left(n \mathfrak{A}-a_{1}-\ldots-a_{n-1}\right) .
$$

Then $\phi$ is continuous in the closed domain

$$
a_{1} \geqq 0, \ldots, \quad a_{n-1} \geqq 0, \quad a_{1}+a_{2}+\ldots+a_{n-1} \leqq n \mathfrak{A} .
$$

It therefore attains a maximum for some system of values $a_{1}^{*}, \ldots, a_{n-1}^{*}$ in the domain.

This proof, the most familiar of all proofs of the theory, is due (so far as we have been able to trace it) to Maclaurin (2). Maclaurin states the theorem in geometrical language, as follows: 'If the Line $A B$ is divided into any Number of Parts $A C, C D$, $D E, E B$, the Product of all those Parts multiplied into one another will be a Maximum when the Parts are equal amongst themselves'. His proof is substantially that which follows. The proof has been rediscovered or reproduced by many later writers, for example by Grebe (1), Chrystal (1, 47).
Cauchy's proof ( $\$ 2.5$ ) may be regarded as a more sophisticated form of Maclaurin's, since he proves the theorem in the special case when $n=2^{m}$ by a process similar to Maclaurin's. In general, Maclaurin's proof is not a 'înite' proof. As we have stated it, it depends on Weierstrass's theorem on the maximum of a continuous function. This would naturally have been taken for granted by Maclaurin (and has also been taken for granted by many of his modern followers, such as Grebe and Chrystal).

It is possible to avoid an appeal to Weierstrass's theorem, but at considerable cost. It is plain that if $a_{1}{ }^{1}, a_{2}{ }^{1} ; a_{1}{ }^{2}, a_{2}{ }^{2} ; \ldots$ are the smallest and largest of the sets resulting from 1,2, $\ldots$ repetitions of Maclaurin's process, then ${a_{1}}^{s}$ increases and $a_{2}{ }^{s}$ decreases as $s$ increases, so that

$$
a_{1}^{s} \rightarrow \alpha_{1}, \quad a_{2}^{s} \rightarrow \alpha_{2}, \quad \alpha_{2} \geqq \alpha_{1} .
$$

A little consideration will show that $n$ repetitions of the process diminish the greatest difference of the $a$ by at least one-half, so that $a_{2}{ }^{n}-a_{1}{ }^{n} \leqq \frac{1}{2}\left(a_{2}-a_{1}\right)$. Hence $a_{2}{ }^{s}-a_{1}{ }^{s} \rightarrow 0$, and $\alpha_{1}=\alpha_{2}$. It follows that all the $a$ tend to the same limit $\mathfrak{A}$. This gives a proof of the theorem, but one a good deal less simple than that in the text.

The reader should work out the analogous proof in which ( $F$ is kept constant and $a_{1}$ and $a_{2}$ are each replaced by $\sqrt{ }\left(a_{1} a_{2}\right)$.
(ii) There is a variation of Cauchy's proof which illustrates a point of some logical importance.

An ordinary inductive proof proceeds from $n$ to $n+1$; the truth of a proposition $P(n)$ follows from the hypotheses
(a) $P(n)$ implies $P(n+1)$,
(b) $P(n)$ is true for $n=1$.

There is another mode of proof which may be called proof by 'backward induction'; the truth of $P(n)$ follows from
( $a^{\prime}$ ) $P(n)$ implies $P(n-1)$,
( $b^{\prime}$ ) $P(n)$ is true for an infinity of $n$.
Cauchy's proof may be arranged as a proof of this last type. First, Cauchy proves ( $b^{\prime}$ ) for $n=2^{m}$. Next, if the theorem is true for $n$, and if $\mathfrak{A}$ is the arithmetic mean of $a_{1}, a_{2}, \ldots, a_{n-1}$, then an application of the theorem to the $n$ numbers $a_{1}, \ldots, a_{n-1}, \mathfrak{A}$ gives

$$
\mathfrak{A}^{n}=\left(\frac{a_{1}+\ldots+a_{n-1}+\mathfrak{H}}{n}\right)^{n}>a_{1} a_{2} \ldots a_{n-1} \mathfrak{A}
$$

the result for $n-1$.
(iii)a Defining $a_{1}$ and $a_{2}$ as in (i), we may replace $a_{1}$ and $a_{2}$ by $\mathfrak{A}$ and $a_{1}+a_{2}-\mathfrak{A}$. Then $\mathfrak{A}$ is again unchanged, and

$$
\mathfrak{U}\left(a_{1}+a_{2}-\mathfrak{U}\right)-a_{1} a_{2}=\left(\mathfrak{A}\left(-a_{1}\right)\left(a_{2}-\mathfrak{U}\right)>0\right.
$$

so that $\mathscr{S}$ is increased. Repeating the process we arrive, after at most $n-1$ steps, at a system of $a$ all equal to $\mathfrak{N}$. It follows that (3) $<\mathfrak{A}$.

This proof is a little more sophisticated but entirely elementary. There is an alternative, which we leave to the reader, in which $a_{1}$ and $a_{2}$ are replaced by $\mathbb{E S}$ and $a_{1} a_{2} / \mathscr{S}$.
(iv) There are a number of inductive proofs of the theorem: see, for example, Chrystal (1,46), Muirhead (3). One of the simplest runs as follows ${ }^{\mathrm{b}}$. Suppose that $0<a_{1} \leqq a_{2} \leqq \ldots \leqq a_{n}, a_{1}<a_{n}$, that $\mathfrak{A}_{\nu}$ and $\mathfrak{G}_{\nu}$ refer to the first $\nu$ of the $a$, and that it has been proved that $\mathfrak{A}_{n-1} \geqq \mathfrak{G}_{n-1}$. Then $a_{n}>\mathfrak{A}_{n-1}$, by Theorem 1, and

$$
\mathfrak{A}_{n}=\frac{(n-1) \mathfrak{A}_{n-1}+a_{n}}{n}=\mathfrak{A}_{n-1}+\frac{a_{n}-\mathfrak{A}_{n-1}}{n}
$$

a For these proofs see Sturm (1, 3), Crawford (1), Briggs and Bryan (1, 185), Muirhead (3), Hardy (1, 32).
${ }^{\text {b }}$ Another simple proof due to R. Rado is given at the end of the chapter (Theorem 60).

Raising this equation to the $n$th power, and remembering that $n>1$, we obtain

$$
\mathfrak{A}_{n}^{n}>\mathfrak{A}_{n-1}^{n}+n \mathfrak{A}_{n-1}^{n-1} \frac{a_{n}-\mathfrak{A}_{n-1}}{n}=a_{n} \mathfrak{A}_{n-1}^{n-1} \geqq a_{n} \mathfrak{F}_{n-1}^{n-1}=\mathfrak{G}_{n}^{n} .
$$

(v) Another interesting proof was given very recently by Steffensen (1,2). It starts from the lemma: if $a_{\nu-1} \leqq a_{\nu}, b_{\nu-1} \leqq b_{\nu}$, and $a_{\nu} \leqq b_{\nu}$, for all $\nu$, then $\Sigma a \Sigma b$ is not decreased by exchanging $a_{i}$ and $b_{i}$, and is increased except when $a_{i}=b_{i}$ or $a_{\nu}=b_{\nu}$ for $\nu \neq i$. The lemma follows at once from the identity

$$
\begin{aligned}
\left\{\Sigma a+\left(b_{i}-a_{i}\right)\right\}\{\Sigma b+ & \left.\left(a_{i}-b_{i}\right)\right\} \\
& =\Sigma a \Sigma b+\left(b_{i}-a_{i}\right)\left\{\left(\Sigma b-b_{i}\right)-\left(\Sigma a-a_{i}\right)\right\} .
\end{aligned}
$$

To deduce the theorem of the means, we write it in the form

$$
\left(a_{1}+a_{1}+\ldots+a_{1}\right) \ldots\left(a_{n}+a_{n}+\ldots+a_{n}\right) \leqq\left(a_{1}+a_{2}+\ldots+a_{n}\right) \ldots
$$

$$
\left(a_{1}+a_{2}+\ldots+a_{n}\right) .
$$

If we suppose, as we may, that $a_{1} \leqq a_{2} \leqq \ldots \leqq a_{n}$, and exchange $n-1$ terms of the first factor of the left-hand side against one term of each of the other factors, we obtain

$$
\left(a_{1}+a_{2}+a_{3}+\ldots+a_{n}\right)\left(a_{1}+a_{2}+a_{2}+\ldots+a_{2}\right) \ldots\left(a_{1}+a_{n}+a_{n}+\ldots+a_{n}\right),
$$

which is greater, by the lemma, unless all the $a$ are equal. The theorem follows by repetition of the argument.
(vi) Further proofs of Theorem 9 (or of the special case considered in this section) are given in $\S \S 2.14,2.21,3.11$, and 4.2 .

### 2.7. Hölder's inequality and its extensions. Ournextgroup

 of theorems centres round Theorem 11 (Hölder's inequality) ${ }^{\text {a }}$.10. Suppose that $(a),(b), \ldots,(l)$ are $m$ sets each of $n$ numbers. Then

$$
\begin{equation*}
\mathfrak{G}(a)+\mathscr{S}(b)+\ldots+\mathfrak{G}(l)<\mathfrak{G}(a+b+\ldots+l), \tag{2.7.1}
\end{equation*}
$$

unless either (1) every two of (a), (b), ..., (l) are proportional, or (2) there is a $\nu$ such that $a_{\nu}=b_{\nu}=\ldots=l_{\nu}=0$.

The theorem states that, if $\Sigma q=1$, then

$$
\begin{aligned}
& a_{1}{ }^{q_{1}} a_{2}^{q_{2}} \ldots a_{n}^{q_{n}}+b_{1}^{q_{1}} b_{2} q_{2} \ldots b_{n}^{q_{n}}+\ldots+l_{1} q_{1} l_{2} q_{2} \ldots l_{n}^{q_{n}} \\
&<\left(a_{1}+b_{1}+\ldots+l_{1}\right)^{q_{1}}\left(a_{2}+b_{2}+\ldots+l_{2}\right)^{q_{2}} \ldots,
\end{aligned}
$$

unless every two columns of the array

$$
\begin{array}{llll}
a_{1}, & b_{1}, & \ldots, & l_{1} \\
a_{2}, & b_{2}, & \ldots, & l_{2} \\
\ldots, & \ldots, & \ldots, & \ldots
\end{array}
$$

a Strictly, 'Hölder's inequality' is Theorem 14, or (2.8.3) of Theorem 13. The inequality (2.7.1) was stated explicitly, for two sets and equal weights, by Minkowski (1, 117).
are proportional or there is a row containing only zeros. A necessary and sufficient condition that all columns should be proportional (i.e. that every pair of columns should be proportional) is that $a_{\mu} b_{\nu}-a_{\nu} b_{\mu}=0, a_{\mu} c_{\nu}-a_{\nu} c_{\mu}=0, \ldots$, for every $\mu$ and $\nu$; and this condition is also necessary and sufficient for the proportionality of all rows. If we remember this, change our notation as between rows and columns of the array, and write $\alpha, \beta, \ldots, \lambda$ for $q_{1}, q_{2}, \ldots, q_{n}$, we see that Theorem 10 is equivalent to
11. If $\alpha, \beta, \ldots, \lambda$ are positive and $\alpha+\beta+\ldots+\lambda=1$, then
$\Sigma a^{\alpha} b^{\beta} \ldots l^{\lambda}<(\Sigma a)^{\alpha}(\Sigma b)^{\beta} \ldots(\Sigma l)^{\lambda}$,
unless either (1) the sets (a), (b), $\ldots$, ( $l$ ) are all proportional, or (2) one set is nul.

The conditions for equality might also be expressed by saying that there is one set which is proportional to all the others (the nul set being proportional to all other sets). The case in which one set is nul is trivial, and we may ignore it in the proof.

Here again we give two proofs.
(i) By Theorem 7,
$(\Sigma a b)^{2}<\Sigma a^{2} \Sigma b^{2}$
unless ( $a$ ) and ( $b$ ) are proportional. Hence

$$
(\Sigma a b c d)^{4} \leqq\left(\Sigma a^{2} b^{2}\right)^{2} \Sigma\left(c^{2} d^{2}\right)^{2} \leqq \Sigma a^{4} \Sigma b^{4} \Sigma c^{4} \Sigma d^{4},
$$

with inequality somewhere unless $(a),(b),(c),(d)$ are proportionala . Repeating the argument we see that

$$
\begin{equation*}
(\Sigma a b \ldots l)^{2^{m}}<\Sigma a^{2 m} \Sigma b^{2^{m}} \ldots \Sigma l^{2^{m}}, \tag{2.7.3}
\end{equation*}
$$

with $2^{m}$ sets $(a),(b), \ldots$, unless all the sets are proportional. This is equivalent to (2.7.2) when every index is $2^{-m}$.

Suppose next that $M$ is any number less than $2^{m}$, and let $(g)$ be the $M$ th set. If ( $a b \ldots g$ ) is not nul, we define $A, B, \ldots, L$ by

$$
\begin{array}{cl}
A^{2^{m}}=a^{M}, \ldots, G^{2^{m}}=g^{M} & (M \text { sets }), \\
H^{2^{m}}=K^{2^{m}}=\ldots=L^{2^{m}}=a b \ldots g & \left(2^{m}-M \text { sets }\right),
\end{array}
$$

so that $A B \ldots L=a b \ldots g$, and apply (2.7.3) to $A, B, \ldots, L$. We thus obtain or

$$
\begin{gather*}
\text { (2.7.4) obtain }(\Sigma a b \ldots g)^{2^{m}}<\Sigma a^{M} \ldots \Sigma g^{M}(\Sigma a b \ldots g)^{2^{m}-M} \\
(\Sigma a b \ldots g)^{M}<\Sigma a^{M} \Sigma b^{M} \ldots \Sigma g^{M}, \tag{2.7.4}
\end{gather*}
$$

a The nul set being excluded, proportionality is now transitive: see $\S 2.1$.
unless the sets $(A),(B), \ldots,(L)$, and so the sets $(a),(b), \ldots,(g)$, are proportional. This is equivalent to (2.7.2) with every index $1 / M$. We have supposed ( $a b \ldots g$ ) not nul; if it is nul then (2.7.4) is obviously true, since none of $(a),(b), \ldots,(g)$ is nul.

If now $\alpha, \beta, \ldots$ are rational, we can write

$$
\alpha=\frac{\alpha^{\prime}}{M}, \quad \beta=\frac{\beta^{\prime}}{M}, \quad \ldots
$$

where $\alpha^{\prime}, \beta^{\prime}, \ldots$ are integers and $\Sigma \alpha^{\prime}=M$. Applying (2.7.2), with every index $1 / M$, to $M$ sets formed by $\alpha^{\prime}$ like sets of $a, \beta^{\prime}$ sets of $b$, and so on, we obtain (2.7.2) with indices $\alpha, \beta, \ldots$.

Finally, when $\alpha, \beta, \ldots$ are not all rational, we replace them by rational approximations whose sum is 1 , form (2.7.2) for these rational indices, and proceed to the limit. In this process ' $<$ ' degenerates into ' $\leqq$ ' and, as in § 2.5 (i), we do not at first obtain a complete proof. We can complete the proof as follows. We can write $\alpha=\alpha_{1}+\alpha_{2}, \beta=\beta_{1}+\beta_{2}, \ldots$, where all the numbers are positive and those with suffix 1 are rational. If then $\Sigma \alpha_{1}=\sigma_{1}$, $\Sigma \alpha_{2}=\sigma_{2}$, so that $\sigma_{1}+\sigma_{2}=1$, and $P_{1} \sigma_{1}=a^{\alpha_{1}} b^{\beta_{1}} \ldots, P_{2}{ }^{\sigma_{2}}=a^{\alpha_{2}} b^{\beta_{2}} \ldots$, we have $\quad \Sigma a^{\alpha} b^{\beta} \ldots l^{\lambda}=\Sigma P_{1} \sigma_{1} P_{2} \sigma_{2} \leqq\left(\Sigma P_{1}\right)^{\sigma_{1}}\left(\Sigma P_{2}\right)^{\sigma_{2}}$.
Since $\alpha_{1}, \beta_{1}, \ldots$ are rational

$$
\Sigma P_{1}=\Sigma a^{\alpha_{1} / \sigma_{1}} \ldots l^{\lambda_{1} / \sigma_{1}}<(\Sigma a)^{\alpha_{1} / \sigma_{1}} \ldots(\Sigma l)^{\lambda_{1} / \sigma_{1}}
$$

while for $\Sigma P_{2}$ we have a similar inequality, but with ' $\leqq$ ' only. Combining our results we obtain (2.7.2).
(ii) We may deduce Theorem 11 from Theorem 9. We have in fact (since no set is nul)

$$
\begin{aligned}
& \frac{\Sigma a^{\alpha} b^{\beta} \ldots l^{\lambda}}{(\Sigma a)^{\alpha}(\Sigma b)^{\beta} \ldots(\Sigma l)^{\lambda}}=\Sigma\left(\frac{a}{\Sigma a}\right)^{\alpha}\left(\frac{b}{\Sigma b}\right)^{\beta} \ldots\left(\frac{l}{\Sigma l}\right)^{\lambda} \\
& \quad \leqq \Sigma\left(\alpha \frac{a}{\Sigma a}+\beta \frac{b}{\Sigma b}+\ldots+\lambda \frac{l}{\Sigma l}\right)=\alpha+\beta+\ldots+\lambda=1
\end{aligned}
$$

There can be equality only if

$$
\frac{a_{\nu}}{\Sigma a}=\frac{b_{\nu}}{\Sigma b}=\ldots=\frac{l_{\nu}}{\Sigma l} \quad(\nu=1,2, \ldots, n)
$$

i.e. if $(a),(b), \ldots,(l)$ are proportional.

It will be observed that, whether $\alpha, \beta, \ldots$ are rational or not, no limiting processes are involved in the proof beyond those already present in the proof of Theorem 9. The principle of the proof is the same as that of the proof of Theorem 13 below given independently by Francis and Littlewood ${ }^{\text {a (1) }}$ ) and F. Riesz (6).
2.8. Hölder's inequality and its extensions (continued). If we suppose $r \neq 0$, and replace $a, b, \ldots, l$ in Theorem 11 by $q a^{r / \alpha}$, $q b^{r / \beta}, \ldots, q l^{r / \lambda}$, we obtain
12. If $r, \alpha, \beta, \ldots, \lambda$ are positive and $\alpha+\beta+\ldots+\lambda=1$, then

$$
\mathfrak{M}_{r}(a b \ldots l)<\mathfrak{M}_{r / \alpha}(a) \mathfrak{M}_{r / \beta}(b) \ldots \mathfrak{M}_{r / \lambda}(l)
$$

unless $\left(a^{1 / \alpha}\right),\left(b^{1 / \beta}\right), \ldots,\left(l^{1 / \lambda}\right)$ are proportional or one of the factors on the right-hand side is zero. If $r<0$, the inequality is reversed.

It is to be observed that, when $r>0$, the second exceptional case occurs only if one of the sets $(a),(b), \ldots$ is nul, whereas when $r<0$ it occurs if any number of any set is zero. When $r=0$ there is equality in any case.

We shall often find it convenient, when we are concerned with two sets of numbers only, to use the notation

$$
\begin{equation*}
k^{\prime}=\frac{k}{k-1} \tag{2.8.1}
\end{equation*}
$$

$k$ being any real number except 1 . The relation (2.8.1) may also be written in the symmetrical forms

$$
\begin{equation*}
(k-1)\left(k^{\prime}-1\right)=1, \quad \frac{1}{k}+\frac{1}{k^{\prime}}=1 \tag{2.8.2}
\end{equation*}
$$

(the last form failing when $k=0, k^{\prime}=0$ ). We say that $k$ and $k^{\prime}$ are conjugate.
13. Suppose that $k \neq 0, k \neq 1$, and that $k^{\prime}$ is conjugate to $k$. Then

$$
\begin{equation*}
\Sigma a b<\left(\Sigma a^{k}\right)^{1 / k}\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}} \quad(k>1) \tag{2.8.3}
\end{equation*}
$$

unless ( $a^{k}$ ) and ( $b^{k^{\prime}}$ ) are proportional; and

$$
\begin{equation*}
\Sigma a b>\left(\Sigma a^{k}\right)^{1 / k}\left(\Sigma b^{k k^{\prime}}\right)^{1 / k^{\prime}} \quad(k<1) \tag{2.8.4}
\end{equation*}
$$

unless either $\left(a^{k}\right)$ and $\left(b^{k^{\prime}}\right)$ are proportional or $(a b)$ is nul.

Cauchy's inequality (Theorem 7) is the special case $k=k^{\prime}=2$, in which $k$ is conjugate to itself.
(i) Suppose that $k>1$. Then (2.8.3) is the special case of Theorem 11 in which there are two sets of letters and $\alpha=1 / k$, $\beta=1 / k^{\prime}$. This is the ordinary form of Hölder's inequality ${ }^{2}$.
(ii) Suppose that $0<k<1$, so that $k^{\prime}<0$. If any $b$ is 0 then the second factor on the right-hand side of (2.8.4) is, as in $\S 2.1$, to be interpreted as 0 , so that (2.8.4) is true unless ( $a b$ ) is nul. If every $b$ is positive, we define $l, u, v$ by
so that

$$
l=1 / k,
$$

and $\quad u=(a b)^{k}, \quad v=b^{-k}$,
so that $\quad a b=u^{l}, \quad a^{k}=u v, \quad b^{k^{\prime}}=v^{l^{\prime}}$.
Then (2.8.4) reduces to (2.8.3) with $u, v, l$ in place of $a, b, k$. The exceptional case is that in which $\left(u^{l}\right)$ and $\left(v^{l}\right)$, i.e. ( $a b$ ) and $\left(b^{k^{\prime}}\right)$, are proportional. If this is so then (since the $b$ are now all positive) the sets ( $a$ ) and ( $b^{k^{\prime}-1}$ ), and therefore the sets ( $a^{k}$ ) and ( $b^{k^{\prime}}$ ), are proportional.
(iii) If $k<0$, then $0<k^{\prime}<1$. This case is reduced to (ii) by exchanging $a$ and $b, k$ and $k^{\prime}$. Both (ii) and (iii) are included in (2.8.4).

The inequalities remain true in the excluded cases $k=0, k=1$ if we adopt appropriate conventions. If $k=0, k^{\prime}=0$, we must interpret (2.8.4) as

$$
a_{1} b_{1}+a_{2} b_{2}+\ldots+a_{n} b_{n}>n\left(a_{1} \ldots a_{n} b_{1} \ldots b_{n}\right)^{1 / n} .
$$

If $k=1$ we may interpret $k^{\prime}$ as $+\infty$ or as $-\infty$. In the first case we interpret (2.8.3) as $\Sigma a b<\operatorname{Max} b \Sigma a$, and in the second we interpret (2.8.4) as $\Sigma a b>\operatorname{Min} b \Sigma a$. We may leave it to the reader to pick out the cases of equality.

We can combine (2.8.3) and (2.8.4) in the single inequality

$$
\begin{equation*}
(\Sigma a b)^{a k^{\prime}}<\left(\Sigma a^{k}\right)^{k^{\prime}}\left(\Sigma b^{k^{\prime}}\right)^{k} \quad(k \neq 0, k \neq 1) . \tag{2.8.5}
\end{equation*}
$$

In view of the extreme importance of Hölder's inequality, we

[^6]depart from our usual practice here and state explicitly the derivative theorem for complex $a, b$.
14. If $k>1$, and $k^{\prime}$ is the conjugate of $k$, then
$$
|\Sigma a b| \leqq\left(\Sigma|a|^{k}\right)^{1 / k}\left(\Sigma|b|^{k^{\prime}}\right)^{1 / k^{\prime}}
$$

There is equality if and only if $\left(\left|a_{\nu}\right|^{k}\right)$ and $\left(\left|b_{\nu}\right|^{k^{\prime}}\right)$ are proportional and $\arg a_{\nu} b_{\nu}$ is independent of $\nu$.

The only additional remark needed for the proof is that

$$
|\Sigma a b|<\Sigma|a b|
$$

unless $\arg a_{\nu} b_{\nu}$ is independent of $\nu$. We regard 0 as having any argument we please.

The following variant of the first part of Theorem 13 is sometimes called 'the converse of Hölder's inequality'.
15. Suppose that $k>1$, that $k^{\prime}$ is conjugate to $k$, and that $B>0$. Then a necessary and sufficient condition that $\Sigma a^{k} \leqq A$ is that $\Sigma a b \leqq A^{1 / k} B^{1 / k^{\prime}}$ for all b for which $\Sigma b^{k^{\prime}} \leqq B$.

The condition is necessary, by (2.8.3). If $\Sigma a^{k}>A$, we can choose the $b$ so that $\Sigma b^{k^{\prime}}=B$ and $\left(b^{k^{\prime}}\right)$ is proportional to $\left(a^{k}\right)$, and then

$$
\Sigma a b=\left(\Sigma a^{k}\right)^{1 / k}\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}}>A^{1 / k} B^{1 / k^{\prime}}
$$

Hence the condition is also sufficient.
Theorem 15 is often useful for the purpose of determining an upper bound for $\Sigma a^{k}$. Any argument based on it can be changed into one which involves only a special (b), but the form stated here, with arbitrary $(b)$, is sometimes more convenienta.
2.9. General properties of the means $\mathfrak{M}_{r}(a)$. We can now prove a theorem which completes and supersedes some of those of §§ 2.3-4.
16.b If $r<s$ then

$$
\begin{equation*}
\mathfrak{M}_{r}(a)<\mathfrak{M}_{s}(a), \tag{2.9.1}
\end{equation*}
$$

unless the $a$ are all equal, or $s \leqq 0$ and an a is zero.
We have proved this already in the special cases (i) $r=-\infty$
${ }^{\text {a }}$ Compare $\S \S 6.9$ (p. 142) and 6.13 (p. 149).
${ }^{\text {b }}$ Schlömilch (1). See also Reynaud and Duhamel (1, 155) and Chrystal (1, 48).
(Theorem 5), (ii) $s=+\infty$ (Theorem 5), (iii) $r=0, s=1$ (Theorem 9), (iv) $s=2 r$ (Theorem 6).

Suppose first that $0<r<s$, and write $r=s \alpha$, so that $0<\alpha<1$, and

$$
p a^{s}=u, \quad p=v
$$

so that $v>0$ and $p a^{s \alpha}=\left(p a^{s}\right)^{\alpha} p^{1-\alpha}=u^{\alpha} v^{1-\alpha}$.
Then

$$
\begin{equation*}
\Sigma u^{\alpha} v^{1-\alpha}<(\Sigma u)^{\alpha}(\Sigma v)^{1-\alpha} \tag{2.9.2}
\end{equation*}
$$

by Theorem 11, unless $u_{\nu} / v_{\nu}$ is independent of $\nu$, i.e. unless $a_{\nu}$ is independent of $\nu$. Hence

$$
\left(\frac{\Sigma p a^{s \alpha}}{\Sigma p}\right)^{1 / s \alpha}<\left(\frac{\Sigma p a^{s}}{\Sigma p}\right)^{1 / s}
$$

which is (2.9.1).
The cases in which $r \leqq 0$ and an $a$ is zero are trivial and we may ignore them. If every $a$ is positive, and $r=0<s$, we have

$$
\left(\mathfrak{M}_{0}(a)\right)^{s}=(\mathfrak{G}(a))^{s}=\mathfrak{F}\left(a^{s}\right)<\mathfrak{H}\left(a^{s}\right)=\left(\mathfrak{M}_{s}(a)\right)^{s},
$$

by Theorem 9 and (2.2.7). The two remaining cases, $r<s<0$ and $r<s=0$, reduce to those already discussed in virtue of (2.2.9).
17. ${ }^{\text {a }}$ If $0<r<s<t$ then

$$
\begin{equation*}
\mathfrak{M}_{s}^{s}<\left(\mathfrak{M}_{r}^{r}\right)^{\frac{t-s}{t-r}}\left(\mathfrak{M}_{t}^{t}\right)^{\frac{s-r}{t-r}} \tag{2.9.3}
\end{equation*}
$$

unless all the a which are not zero are equal.
We restrict the parameters to be positive, the complications introduced by negative or zero values being hardly worth pursuing systematically.

We may write

$$
s=r \alpha+t(1-\alpha) \quad(0<\alpha<1)
$$

The inequality is then

$$
\Sigma q a^{s}<\left(\Sigma q a^{r}\right)^{\alpha}\left(\Sigma q a^{t}\right)^{1-\alpha}
$$

and reduces to a case of Theorem 11 when we write $u=q \alpha^{r}$, $v=q a^{t}$. The condition for equality is that $(u)$ and $(v)$ should be proportional, and this is plainly equivalent to that stated in the enunciation. The reader should observe the difference between the conditions for equality in Theorems 16 and 17.

We shall see later (§3.6, Theorem 87) that Theorem 17 may be stated in a more striking form.
2.10. The sums $\mathbb{S}_{r}(a)$. (i) We write

$$
\mathfrak{S}_{r}=\Im_{r}(a)=\left(\Sigma a^{r}\right)^{1 / r} \quad(r>0)
$$

We confine our attention to positive $r$, leaving the construction of a theory of $\mathbb{S}_{r}$ for $r \leqq 0$ as an exercise to the reader.
18. If $0<r<s<t$ then

$$
\begin{equation*}
\mathfrak{S}_{s}^{s}<\left(\mathfrak{S}_{r}^{r}\right)^{\frac{t-s}{t-r}}\left(\mathbb{S}_{t}^{t}\right)^{\frac{s-r}{t-r}} \tag{2.10.1}
\end{equation*}
$$

unless all the a which are not zero are equal.
This is essentially the same theorem as Theorem 17. In fact

$$
\begin{equation*}
\mathfrak{S}_{r}(a)=n^{1 / r} \mathfrak{M}_{r}(a), \tag{2.10.2}
\end{equation*}
$$

the mean $\mathfrak{M}_{r}(a)$ being formed with unit weights, and (2.10.1) reduces to (2.9.3), the powers of $n$ disappearing.

The correspondence between Theorems 17 and 18 depends essentially on the fact that (2.9.3) and (2.10.1) are homogeneous in the second sense of $\S 1.4$, namely in the sign $\Sigma$. There is a theorem for sums corresponding to Theorem 16, but in this theorem, which is expressed by (2.10.3) below, the sign of inequality is reversed; (2.10.3) is not homogeneous in $\Sigma$, and is not related to (2.9.1) as (2.10.1) is related to (2.9.3).
19. ${ }^{\text {a }}$ If $0<r<s$ then
(2.10.3)

$$
\mathfrak{S}_{s}(a)<\mathfrak{S}_{r}(a)
$$

unless all the a but one are zero.
Since the inequality is homogeneous in the $a$, we may suppose $\Sigma a^{r}=1$, i.e. $\mathfrak{S}_{r}=1$.b Then $a_{\nu} \leqq 1$ for every $\nu$, and so $a_{\nu}{ }^{s} \leqq a_{\nu}^{r}$ and

$$
\Sigma a^{s} \leqq \Sigma a^{r}=1
$$

If more than one $a$ is positive then at least one positive $a$ is less than 1 , and then there is inequality. Theorem 19 is usually quoted as Jensen's inequality.
(ii) We add the theorems for $\mathbb{S}_{r}(a)$ corresponding to Theorems 4 and 3.

[^7]20. $\mathbb{S}_{r} \rightarrow \operatorname{Max} a$ when $r \rightarrow \infty$.
21. $\mathfrak{S}_{r} \rightarrow \infty$ when $r \rightarrow 0$, unless all the a but one are zero.

Theorem 20 follows from (2.10.2) and Theorem 4. To prove Theorem 21 we have only to observe that $\Sigma a^{r}=N+o(1)$, where $N$ is the number of positive $a$.
(iii) Theorem 19, combined with Theorem 11, gives the following theorem of Jensen ${ }^{\text {a }}$.
22. If $\alpha, \beta, \ldots, \lambda$ are positive and $\alpha+\beta+\ldots+\lambda>1$, then

$$
\Sigma a^{\alpha} b^{\beta} \ldots l^{\lambda}<(\Sigma a)^{\alpha}(\Sigma b)^{\beta} \ldots(\Sigma l)^{\lambda},
$$

unless every number of one set or all but one of each set is zero, and, in the latter case, those which are positive have the same rank.

We can write $\alpha=k \alpha^{\prime}, \beta=k \beta^{\prime}, \ldots$, where $k>1$ and $\alpha^{\prime}+\beta^{\prime}+\ldots=1$. If then $a^{k}=A, b^{k}=B, \ldots$, we have

$$
\begin{aligned}
\Sigma a^{\alpha} b^{\beta} \ldots l^{\lambda} & =\Sigma A^{\alpha^{\prime}} B^{\beta^{\prime}} \ldots L^{\lambda^{\prime}} \leqq(\Sigma A)^{\alpha^{\prime}}(\Sigma B)^{\beta^{\prime}} \ldots(\Sigma L)^{\lambda^{\prime}} \\
& =\left(\Sigma a^{k}\right)^{\alpha^{/ k}} \ldots\left(\Sigma l^{k}\right)^{\lambda / k} \leqq(\Sigma a)^{\alpha} \ldots(\Sigma l)^{\lambda},
\end{aligned}
$$

by Theorems 11 and 19. There is inequality somewhere unless the conditions for equality in both theorems are satisfied.
(iv) It is natural to consider weighted sums

$$
\mathfrak{T}_{r}=\mathfrak{T}_{r}(a)=\mathfrak{T}_{r}(a, p)=\left(\Sigma p a^{r}\right)^{1 / r} .
$$

It is plain that there can be no universal relation of the type (2.9.1) or (2.10.3), since $\mathfrak{I}_{r}$ is the $\mathfrak{S}_{r}$ of Theorem 19 when $p_{\nu}=1$ and is $\mathfrak{M}_{r}$ when $\Sigma p_{\nu}=1$. The possibilities in this direction are settled by the following theorem.
23. A necessary and sufficient condition that
(2.10.4)

$$
\mathfrak{T}_{r} \leqq \mathfrak{T}_{s} \quad(0<r<s),
$$

for given weights $p$ and all $a$, is that $\Sigma p \leqq 1$. There is then inequality unless (a) is nul, or $\Sigma p=1$ and all the $a$ are equal.

A necessary and sufficient condition that
(2.10.5)

$$
\mathfrak{I}_{s} \leqq \mathfrak{I}_{r} \quad(0<r<s),
$$

for given weights $p$ and all $a$, is that $p_{\nu} \geqq 1$ for every $\nu$. There is then inequality unless (a) is nul, or $a_{k}>0, p_{k}=1$, and the remaining a are zero.
(i) If we take $a_{\nu}=1$ for every $\nu$, then $\mathfrak{I}_{r}=(\Sigma p)^{1 / r}$, and (2.10.4) can be true only if $\Sigma p \leqq 1$. If this condition is satisfied, and $r=s \alpha$, so that $0<\alpha<1$, we have

$$
\Sigma p a^{r}=\Sigma\left(p a^{s}\right)^{\alpha} p^{1-a} \leqq\left(\Sigma p a^{s}\right)^{\alpha}(\Sigma p)^{1-\alpha} \leqq\left(\Sigma p a^{s}\right)^{\alpha},
$$

which is (2.10.4). The conditions for equality are plainly as stated.
a Jensen (2).
(ii) If we take $a_{k}=1$ and the other $a$ zero, $\mathfrak{T}_{r}=p_{k}{ }^{1 / r}$, and (2.10.5) can be true only if $p_{k} \geqq 1$. If we assume that this condition is satisfied, write $s=r \beta$, so that $\beta>1$, and assume, as we may on grounds of homogeneity, that $\Sigma p a^{r}=1$, then $p a^{r} \leqq 1$ for every $\nu$ and

$$
\Sigma p a^{s}=\Sigma\left(p a^{r}\right)^{\beta} p^{1-\beta} \leqq \Sigma\left(p a^{r}\right)^{\beta} \leqq \Sigma p a^{r}
$$

which is (2.10.5). The conditions for equality are again plainly as stated.
2.11. Minkowski's inequality. Our next theorem is a generalisation of Theorem 10.
24. Suppose that $r$ is finite and not equal to 1 . Then

$$
\begin{equation*}
\mathfrak{M}_{r}(a)+\mathfrak{M}_{r}(b)+\ldots+\mathfrak{M}_{r}(l)>\mathfrak{M}_{r}(a+b+\ldots+l) \quad(r>1), \tag{2.11.1}
\end{equation*}
$$

$$
\begin{equation*}
\mathfrak{M}_{r}(a)+\mathfrak{M}_{r}(b)+\ldots+\mathfrak{M}_{r}(l)<\mathfrak{M}_{r}(a+b+\ldots+l) \quad(r<1), \tag{2.11.2}
\end{equation*}
$$

unless (a), (b), $\ldots,(l)$ are proportional, or $r \leqq 0$ and

$$
a_{\nu}=b_{\nu}=\ldots=l_{\nu}=0
$$

for some $\nu$.
There is equality for any $a, b, \ldots$ when $r=1$. Theorem 10 is the special case $r=0$. The main result remains true (and is trivial) when $r=\infty$ or $r=-\infty$, except that the conditions for equality require a restatement which may be left to the reader.

We take the means with $q$, and write

$$
a+b+\ldots+l=s, \quad \mathfrak{M}_{r}(s)=S
$$

Then

$$
\begin{aligned}
S^{r} & =\Sigma q s^{r}=\Sigma q a s^{r-1}+\Sigma q b s^{r-1}+\ldots+\Sigma q l s^{r-1} \\
& =\Sigma\left(q^{1 / r} a\right)\left(q^{1 / r} s\right)^{r-1}+\ldots+\Sigma\left(q^{1 / r} l\right)\left(q^{1 / r} s\right)^{r-1}
\end{aligned}
$$

Suppose first that $r>1$. Applying (2.8.3) of Theorem 13 to each sum on the right, we obtain

$$
\begin{equation*}
S^{r} \leqq\left(\Sigma q a^{r}\right)^{1 / r}\left(\Sigma q s^{r}\right)^{1 / r^{r}}+\ldots=S^{r-1}\left(\left(\Sigma q a^{r}\right)^{1 / r}+\ldots\right) . \tag{2.11.3}
\end{equation*}
$$

There is equality only if $\left(q a^{r}\right),\left(q b^{r}\right), \ldots$ are all proportional to $\left(q s^{r}\right)$, i.e. if $(a),(b), \ldots$ are proportional. Since $S$ is positive (except in the trivial case when every set is nul), this establishes (2.11.1) a.

Suppose next that $0<r<1$. Unless all the sets $(a),(b), \ldots$ are nul, $s_{\nu}>0$ for some $\nu$. If $s_{\nu}=0$ for any particular $\nu$, then ${ }^{\text {a }}$ This proof is due to F. Riesz (1, 45).
$a_{\nu}=b_{\nu}=\ldots=l_{\nu}=0$, and we may omit that value of $\nu$ from consideration. We may therefore argue as if $s_{\nu}>0$ for every $\nu$. In that case (2.8.4) of Theorem 13 gives (2.11.3) with the sign of inequality reversed, and the proof may be completed as before.

Finally, suppose that $r<0$. If any $s_{\nu}$ is zero, all the means are zero; we may therefore assume that $s_{\nu}>0$ for every $\nu$. If any $a_{\nu}$ is zero, $\mathfrak{M l}_{r}(a)=0$, and we may omit the letter $a$. ${ }^{\text {a }}$ We may therefore argue on the assumption that every $a, b, \ldots$ is positive, and then again everything follows from (2.8.4) of Theorem 13.

When the $q$ are equal, we obtain
25. If $r$ is finite and not equal to 0 or 1 , then

$$
\begin{equation*}
\left(\Sigma(a+b+\ldots+l)^{r}\right)^{1 / r}<\left(\Sigma a^{r}\right)^{1 / r}+\ldots+\left(\Sigma l^{r}\right)^{1 / r} \quad(r>1) \tag{2.11.4}
\end{equation*}
$$

$$
\begin{equation*}
\left(\Sigma(a+b+\ldots+l)^{r}\right)^{1 / r}>\left(\Sigma a^{r}\right)^{1 / r}+\ldots+\left(\Sigma l^{r}\right)^{1 / r} \quad(r<1) \tag{2.11.5}
\end{equation*}
$$

unless $(a),(b), \ldots,(l)$ are proportional, or $r<0$ and $a_{\nu}, b_{\nu}, \ldots, l_{\nu}$ are all zero for some $\nu$.

It is (2.11.4) which is usually called Minkowski's inequalityb. Theorem 24 is more general than Theorem 25 in appearance only, since it may be deduced from Theorem 25 by writing $p^{1 / r} a, p^{1 / r} b, \ldots$ for $a, b, \ldots$.

Theorem 24 may be given a very elegant symmetrical form ${ }^{c}$.
26. Suppose that $\mathfrak{M}^{(\mu)}$ denotes a mean taken with respect to the suffix $\mu$, with weights $p_{\mu}$, and $\mathfrak{M}^{(\nu)}$ one taken with respect to $\nu$ with weights $q_{\nu} ;{ }^{\mathrm{d}}$ and that $0<r<s<\infty$. Then

$$
\mathfrak{M}_{s}{ }^{(\nu)} \mathfrak{M}_{r}{ }^{(\mu)}\left(a_{\mu \nu}\right)<\mathfrak{M}_{r}{ }^{(\mu)} \mathfrak{M}_{s}{ }^{(\nu)}\left(a_{\mu \nu}\right),
$$

except when $a_{\mu \nu}=b_{\mu} c_{\nu}$.
The result holds generally for all $r, s$ such that $r<s$, except for the specification of the cases of equality.
a Here we use (2.2.15). b Minkowski (1, 115-117).
c Theorem 26 was communicated to us in 1929 by Mr A. E. Ingham. The same formulation of Minkowski's inequality was found independently by Jessen and published in his paper 1. This and his later papers 2 and 3 contain many interesting generalisations: see Theorems 136 and 137.
d We depart here from our usual convention about $q ; \Sigma q$ is not necessarily 1 (though we prove the inequality by transforming it into one in which we may sup. pose $\Sigma q=1$ ).

We prove the theorem for $0<r<s<\infty$, leaving the other cases to the reader. There are various supplementary cases of equality when $r \leqq 0$ or one of $r$ and $s$ is infinite.

Let $s / r=k>1$ and $p_{\mu} a_{\mu \nu}^{r}=A_{\mu \nu}$. Then the inequality to be proved is
or

$$
\left\{\sum_{\nu=1}^{n} q_{\nu}\left(\sum_{\mu=1}^{m} p_{\mu} a_{\mu \nu}^{r}\right)^{s / r}\right\}^{1 / s}<\left\{\sum_{\mu=1}^{m} p_{\mu}\left(\sum_{\nu=1}^{n} q_{\nu} a_{\mu \nu}^{s}\right)^{r / s}\right\}^{1 / r}
$$

$$
\left\{\sum_{\nu=1}^{n} q_{\nu}\left(\sum_{\mu=1}^{m} A_{\mu \nu}\right)^{k}\right\}^{1 / k}<\sum_{\mu=1}^{m}\left(\sum_{\nu=1}^{n} q_{\nu} A_{\mu \nu}^{k}\right)^{1 / k}
$$

This reduces to (2.11.1) when $\Sigma q=1$, and, being homogeneous in the $q$, is true without this restriction.
2.12. A companion to Minkowski's inequality. The theorem which follows is an analogue of Theorem 25 of a simpler kind.
27. If $r$ is positive and not equal to 1 then

$$
\begin{equation*}
\Sigma(a+b+\ldots+l)^{r}>\Sigma a^{r}+\Sigma b^{r}+\ldots+\Sigma l^{r} \quad(r>1) \tag{2.12.1}
\end{equation*}
$$

(2.12.2) $\quad \Sigma(a+b+\ldots+l)^{r}<\Sigma a^{r}+\Sigma b^{r}+\ldots+\Sigma l^{r} \quad(0<r<1)$,
unless all numbers but one of each set $a_{\nu}, b_{\nu}, \ldots, l_{\nu}(\nu=1,2, \ldots, n)$ are zero.

This follows at once from Theorem 19, since for example

$$
(a+b+\ldots+l)^{r}>a^{r}+b^{r}+\ldots+l^{r}
$$

if $r>1$, unless all of $a, b, \ldots, l$ but one are zero. It should be noticed that the sense of (2.12.1) and (2.12.2) is opposite to that of (2.11.4) and (2.11.5).

What is usually required in practiceis a combination of (2.11.4) and (2.12.2), viz.
28. If $r>0$ then

$$
\left(\Sigma(a+b+\ldots+l)^{r}\right)^{R} \leqq\left(\Sigma a^{r}\right)^{R}+\left(\Sigma b^{r}\right)^{R}+\ldots+\left(\Sigma l^{r}\right)^{R}
$$

where $R=1$ if $0<r \leqq 1$ and $R=1 / r$ if $r>1$.
2.13. Illustrations and applications of the fundamental inequalities. (i) Geometrical interpretations of Hölder's and Min-
kowski's inequalities. Two particularly simple cases of Hölder's and Minkowski's inequalities are

$$
\begin{align*}
& \left(x_{1} x_{2}+y_{1} y_{2}+z_{1} z_{2}\right)^{2}<\left(x_{1}^{2}+y_{1}^{2}+z_{1}^{2}\right)\left(x_{2}^{2}+y_{2}^{2}+z_{2}^{2}\right)  \tag{2.13.1}\\
& \sqrt{ }\left\{\left(x_{1}+x_{2}\right)^{2}+\left(y_{1}+y_{2}\right)^{2}+\left(z_{1}+z_{2}\right)^{2}\right\} \\
& \quad<\sqrt{ }\left(x_{1}^{2}+y_{1}^{2}+z_{1}^{2}\right)+\sqrt{ }\left(x_{2}^{2}+y_{2}^{2}+z_{2}^{2}\right)
\end{align*}
$$

These hold for all real values of the variables, and express the facts that (1) the cosine of a real angle is numerically less than 1 , and (2) the sum of two sides of a triangle is greater than the third side. The exceptional cases are those in which (1) the vectors $\left(x_{1}, y_{1}, z_{1}\right)$ and ( $x_{2}, y_{2}, z_{2}$ ) are parallel (with the same or opposite senses), and (2) the vectors are parallel and have the same sense.

The ordinary form of Minkowski's inequality is the extension of (2.13.2) to space of $n$ dimensions with a generalised definition of distance, viz.

$$
P_{1} P_{2}=\left(\left|x_{1}-x_{2}\right|^{r}+\left|y_{1}-y_{2}\right|^{r}+\ldots\right)^{1 / r} \quad(r \geqq 1)
$$

The most obvious extensions of (2.13.1) are connected not with Hölder's inequality for general $r$ but with a generalisation of the case $r=2$ in a different direction.
29. If $\Sigma a_{\mu \nu} x_{\mu} x_{\nu}$, where $a_{\mu \nu}=a_{\nu \mu}$, is a positive quadratic form (with real, but not necessarily positive, coefficients), then

$$
\left(\Sigma a_{\mu \nu} x_{\mu} y_{\nu}\right)^{2}<\Sigma a_{\mu \nu} x_{\mu} x_{\nu} \Sigma a_{\mu \nu} y_{\mu} y_{\nu}
$$

unless $(x)$ and ( $y$ ) are proportional.
This is an immediate consequence of the fact that

$$
\Sigma a_{\mu \nu}\left(\lambda x_{\mu}+\mu y_{\mu}\right)\left(\lambda x_{\nu}+\mu y_{\nu}\right)
$$

is positive: compare the second proof of Theorem 7. It represents geometrically an extension of (2.13.1) to $n$-dimensional space, with oblique coordinates or a non-Euclidean metric.

To illustrate Theorem 15, take $k=2, A=l^{2}, B=1$, and rectangular coordinates. The theorem then asserts that, if the length of the projection of a vector along an arbitrary direction does not exceed $l$, the length of the vector does not exceed $l$.
(ii) $A$ theorem of Hadamard ${ }^{\text {a }}$. In our next theorem also we are concerned with a set of numbers $a_{\mu \nu}$ real but not necessarily positive.
30. If $D$ is the determinant whose constituents are
then

$$
a_{\mu \nu} \quad(\mu, \nu=1,2, \ldots, n)
$$

$$
\begin{equation*}
D^{2} \leqq \Sigma a_{1 \kappa}^{2} \Sigma a_{2 \kappa}^{2} \ldots \Sigma a_{n \kappa}^{2} \tag{2.13.3}
\end{equation*}
$$

There is equality only when

$$
\begin{equation*}
a_{\mu 1} a_{\nu 1}+a_{\mu 2} a_{\nu 2}+\ldots+a_{\mu n} a_{\nu n}=0 \tag{2.13.4}
\end{equation*}
$$

for every distinct pair $\mu, \nu$, or when one of the factors on the righthand side of (2.13.3) vanishes.

The geometrical significance of the theorem is that the volume of a parallelepiped in $n$-space does not exceed the product of the edges diverging from one corner, and that there is equality only when they are orthogonal or an edge vanishes.

Suppose that $\Sigma c_{\mu \nu} x_{\mu} x_{\nu}$, where $c_{\mu \nu}=c_{\nu \mu}$, is a positive quadratic form, and that $\Delta$ is the determinant whose constituents are $c_{\mu \nu}$. Then the equation

$$
\left|\begin{array}{ccc}
c_{11}-\lambda & c_{12} & \cdots  \tag{2.13.5}\\
c_{21} & c_{22}-\lambda & \ldots \\
\ldots & \ldots & \cdots
\end{array}\right|=0
$$

has $n$ positive roots ${ }^{b}$ whose sum is $\Sigma c_{\mu \mu}$ and whose product is $\Delta$. Hence, by Theorem 9,

$$
\begin{equation*}
\Delta \leqq\left(\frac{c_{11}+c_{22}+\ldots+c_{n n}}{n}\right)^{n} \tag{2.13.6}
\end{equation*}
$$

If $c_{\mu \mu}>0$ for all $\mu$, then the form

$$
\Sigma \frac{c_{\mu \nu}}{\sqrt{ }\left(c_{\mu \mu} c_{\nu \nu}\right)} x_{\mu} x_{\nu}=\Sigma C_{\mu \nu} x_{\mu} x_{\nu}
$$

is also positive; and if we apply (2.13.6) to this form, we obtain

$$
\begin{equation*}
\Delta \leqq c_{11} c_{22} \ldots c_{n n} \tag{2.13.7}
\end{equation*}
$$

[^8]This is substantially equivalent to Hadamard's theorem. For the form

$$
\Sigma\left(a_{1 \kappa} x_{1}+a_{2 \kappa} x_{2}+\ldots+a_{n \kappa} x_{n}\right)^{2}=\Sigma c_{\mu \nu} x_{\mu} x_{\nu}
$$

is positive unless $D=0$. Also $\Delta=D^{2}$ and

$$
c_{\mu \mu}=a_{\mu 1}{ }^{2}+a_{\mu 2}{ }^{2}+\ldots+a_{\mu n}{ }^{2},
$$

so that (2.13.7) is (2.13.3).
For equality in (2.13.6), all the roots of (2.13.5) must be equal, which is only possible if $c_{\mu \nu}=0$ whenever $\mu \neq \nu$ and $c_{\mu \mu}$ is independent of $\mu$. Hence, for equality in (2.13.7), we must have $C_{\mu \nu}=0$ for $\mu \neq \nu, C_{\mu \mu}$ independent of $\mu$. The last condition is certainly satisfied, since $C_{\mu \mu}=1$, and $C_{\mu \nu}=0$ is $c_{\mu \nu}=0$, which is (2.13.4).

We can extend the theorem to determinants with complex constituents by using Hermitian instead of quadratic forms. Further extensions have been made by Schur (2) ${ }^{\text {a }}$.

The following ingenious proof of (2.13.7) is due to Oppenheimb ${ }^{\text {b }}$. Oppenheim's argument establishes not only (2.13.7), and so Hadamard's theorem, but also the inequalities (2.13.8) and (2.13.9) below, due to Minkowskic and Fischerd respectively.

Any two positive quadratic forms $\Sigma c_{i k} x_{i} x_{k}, \Sigma d_{i k} x_{i} x_{k}$ may be reduced simultaneously, by a linear transformation of determinant unity, to sums of squarese, say $\Sigma c_{\nu} y_{\nu}{ }^{2}, \Sigma d_{\nu} y_{\nu}{ }^{2}$, where $c_{\nu}$ and $d_{\nu}$ are positive. Then $\Sigma\left(c_{i k}+d_{i k}\right) x_{i} x_{k}$ is reduced to $\Sigma\left(c_{\nu}+d_{\nu}\right) y_{\nu}{ }^{2}$, and the determinants $\left|c_{i k}\right|, \ldots$ of the forms satisfy

$$
\left|c_{i k}\right|=\Pi c_{\nu}, \quad\left|d_{i k}\right|=\Pi d_{\nu}, \quad\left|c_{i k}+d_{i k}\right|=\Pi\left(c_{\nu}+d_{\nu}\right) .
$$

Hence, applying Theorem 10 to the sets $\left(c_{\nu}\right),\left(d_{\nu}\right)$, we obtain

$$
\begin{equation*}
\left|c_{2 k}\right|^{1 / n}+\left|d_{i k}\right|^{1 / n} \leqq\left|c_{i k}+d_{i k}\right|^{1 / n} . \tag{2.13.8}
\end{equation*}
$$

Suppose now that the matrix of the $d$ is formed from that of the $c$ by multiplying, first the first $r$ rows, and then the first $r$

[^9]columns, by $-1 .{ }^{\text {a }}$ If then we divide (2.13.8) by 2 , and raise to the $n$th power, we obtain
\[

$$
\begin{equation*}
\left|c_{i k}\right|=\left|c_{11} \ldots c_{n n}\right| \leqq\left|c_{11} \ldots c_{r r}\right|\left|c_{r+1, r+1} \ldots c_{n n}\right| \tag{2.13.9}
\end{equation*}
$$

\]

where $\left|c_{11} \ldots c_{r r}\right|$ denotes the north-west diagonal minor of $r$ rows and columns in $\left|c_{i k}\right|$, and $\left|c_{r+1, r+1} \ldots c_{n n}\right|$ denotes the complementary south-east minor. Repeating the argument, replacing each of the factors on the right-hand side of (2.13.9) by two factors, and so on, we ultimately obtain (2.13.7).
(iii) The modulus of a matrix. Suppose that $A$ and $B$ are the matrices of $n$ rows and columns whose elements are $a_{\mu \nu}$ and $b_{\mu \nu}$; the elements may be complex. The matrices $A+B$ and $B A$ are defined as the matrices whose elements are

$$
a_{\mu \nu}+b_{\mu \nu}, \quad b_{\mu 1} a_{1 \nu}+b_{\mu 2} a_{2 \nu}+\ldots+b_{\mu n} a_{n \nu}
$$

31.b If $|A|$, the modulus of the matrix $A$, is defined by

$$
|A|=\sqrt{ }\left\{\Sigma\left|a_{\mu \nu}\right|^{2}\right\},
$$

then

$$
|A+B| \leqq|A|+|B|, \quad|B A| \leqq|B||A| .
$$

The first inequality is an immediate consequence of Theorem 25, with $r=2$. The second follows from Theorem 7, since

$$
\underset{\mu, \nu}{\Sigma\left|b_{\mu 1} a_{1 \nu}+\ldots+b_{\mu n} a_{n \nu}\right|^{2} \leqq \sum_{\mu, \nu, p, q}\left|b_{\mu p}\right|^{2}\left|a_{q \nu}\right|^{2} . . . . ~}
$$

(iv) Maxima and minima in elementary geometry. We quote (as exercises for the reader) a few of the numerous applications of the fundamental inequalities to problems of elementary geometry.
32. The area of a triangle of given perimeter $2 p$ is a maximum if the sides $a, b, c$ are equal.
[Apply Theorem 9 to $p-a, p-b, p-c$.]
33. If the surface of a rectangular parallelepiped is given, the volume is greatest when the parallelepiped is a cube.
[Denote the edges diverging from a corner by $a, b, c$ and apply Theorem 9 to $b c, c a, a b$. There is an analogous theorem for a parallelepiped in $n$ dimensions; if $k<n$, and the surface of the $k$-dimensional boundary is given, the volume is greatest when the parallelepiped is rectangular and its edges are equal. This may be proved by combining Theorems 9 and 30 with identities between determinants.]

[^10]34. Définition. Si la Base d'une Pyramide est circonscriptible à un Cercle; et si le Pié de la Hauteur est au Centre de ce Cercle: J'appelle cette Pyramide droite.

Dans une Pyramide droite toutes les Faces ont une même Hauteur, et sont également inclinées au Plan de la Base.

Théorème. Soient deux Pyramides de même Hauteur, dont les Bases sont égales tant en Surface qu'en Contour; que l'une soit droite et que l'autre ne le soit pas: $j$ 'affirme que la Surface de la première Pyramide est plus petite que la Surface de la seconde.
[Lhuilier (1, 116).]
[Let $h$ be the height, $b_{\nu}$ a side of the base, and $p_{\nu}$ the perpendicular from the foot of the altitude on to $b_{\nu}$. Then the lateral surface of the second pyramid is

$$
\frac{1}{2} \Sigma b_{\nu} \sqrt{ }\left(h^{2}+p_{\nu}{ }^{2}\right)>\frac{1}{2} \sqrt{ }\left\{\left(\Sigma h b_{\nu}\right)^{2}+\left(\Sigma p_{\nu} b_{\nu}\right)^{2}\right\},
$$

by (2.11.4) of Theorem 25, unless all the $p_{\nu}$ are equal.]
(v) Some inequalities useful in elementary analysis. The following theorems, which are easy deductions from Theorem 9, are fundamental in the theory of the exponential and logarithmic functions.
35. If $\xi>0,0<m<n$, then

$$
\left(1+\frac{\xi}{m}\right)^{m}<\left(1+\frac{\xi}{n}\right)^{n} .
$$

If also $\xi<m$, then

$$
\left(1-\frac{\xi}{m}\right)^{-m}>\left(1-\frac{\xi}{n}\right)^{-n}
$$

36. If $\xi>0, \xi \neq 1,0<m<n$, then

$$
n\left(\xi^{1 / n}-1\right)<m\left(\xi^{1 / m}-1\right) .
$$

We have, by Theorem 9,

$$
\left(1+\frac{\xi}{m}\right)^{\frac{m}{n}} 1^{\frac{n-m}{n}}<\frac{m}{n}\left(1+\frac{\xi}{m}\right)+\frac{n-m}{n} 1=1+\frac{\xi}{n}
$$

If $\xi<m$, we may write $-\xi$ for $\xi$. This proves Theorem 35. Theorem 36 follows from Theorem 35 if we replace $\xi$ in Theorem 36 by $\left(1 \pm \frac{\xi}{m}\right)^{m}$.
2.14. Inductive proofs of the fundamental inequalities. Our fundamental theorems are Theorems 9,10 (or 11), and 24 (or 25), which we refer to shortly as G, H, M. We deduced H from $\mathrm{G}^{\mathrm{a}}$ and M from H ; G is a limiting case of $\mathrm{H}, \mathrm{H}$ a special case, or anticipation, of M.

The simplest case of $G$ is
37. $\left(\mathrm{G}_{0}\right): \quad a^{\alpha} b^{\beta}<a \alpha+b \beta \quad(\alpha+\beta=1)$.

We show first that $G$ can be deduced from $G_{0}$ by induction.
a. Though giving also an independent proof of H .

Suppose that G has been proved for $m$ letters $a, b, \ldots, k$ (or for any smaller number), and that

$$
\alpha+\beta+\ldots+\kappa+\lambda=1, \quad \alpha+\beta+\ldots+\kappa=\sigma .
$$

Then

$$
\begin{aligned}
a^{\alpha} b^{\beta} \ldots k^{\kappa} l^{\lambda} & =\left(a^{\alpha / \sigma} b^{\beta / \sigma} \ldots k^{\kappa / \sigma}\right)^{\sigma} l^{\lambda} \\
& \leqq\left(a^{\alpha / \sigma} b^{\beta / \sigma} \ldots k^{k / \sigma}\right) \sigma+l \lambda \leqq a \alpha+b \beta+\ldots+k \kappa+l \lambda,
\end{aligned}
$$

by G for 2 and for $m$ letters. There is equality in the final result only if

$$
a^{\alpha / \sigma} b^{\beta / \sigma} \ldots k^{\kappa / \sigma}=l, \quad a=b=\ldots=k
$$

i.e. if all letters are equal. Hence $\mathbf{G}$ is true for $m+1$ letters.

The simplest cases of H and M are
38. $\left(\mathrm{H}_{0}\right)$ :

$$
a_{1}{ }^{\alpha} b_{1}{ }^{\beta}+a_{2}{ }^{\alpha} b_{2}{ }^{\beta}<\left(a_{1}+a_{2}\right)^{\alpha}\left(b_{1}+b_{2}\right)^{\beta} \quad(\alpha+\beta=1) .
$$

39. $\left(\mathrm{M}_{0}\right)$ :

$$
\left\{\left(a_{1}+b_{1}\right)^{r}+\left(a_{2}+b_{2}\right)^{r}\right\}^{1 / r}<\left(a_{1}^{r}+a_{2}^{r}\right)^{1 / r}+\left(b_{1}^{r}+b_{2}^{r}\right)^{1 / r} \quad(r>1)
$$

(with a reversed inequality when $r<1$ ). We can deduce $\mathrm{H}_{0}$ from $G_{0}$ and $M_{0}$ from $H_{0}$ by specialising our deductions of $H$ from $G$ and of $M$ from $H$. We can also deduce $H$ and $M$ from $H_{0}$ and $M_{0}$ by induction, but, since these inductive proofs are not essential to our argument, we need only sketch them.
(i) We have

$$
\begin{aligned}
& a_{1}{ }^{\alpha} b_{1}{ }^{\beta}+a_{2}{ }^{\alpha} b_{2}{ }^{\beta}+a_{3}{ }^{\alpha} b_{3}{ }^{\beta} \leqq\left(a_{1}+a_{2}\right)^{\alpha}\left(b_{1}+b_{2}\right)^{\beta}+a_{3}{ }^{\alpha} b_{3}{ }^{\beta} \\
& \leqq\left(a_{1}+a_{2}+a_{3}\right)^{\alpha}\left(b_{1}+b_{2}+b_{3}\right)^{\beta} .
\end{aligned}
$$

The process may be repeated, and there is no difficulty in picking out the cases of equality. We thus obtain (2.8.3) of Theorem 13 ( H for two sets of $n$ numbers).
Next, if $\alpha+\beta+\gamma=1, \alpha+\beta=\sigma$, we have
$\Sigma \alpha^{\alpha} b^{\beta} c^{\gamma}=\Sigma\left(a^{\alpha / \sigma} b \beta / \sigma\right)^{\sigma} c^{\gamma} \leqq\left(\Sigma a^{\alpha / \sigma} b^{\beta / \sigma}\right)^{\sigma}(\Sigma c)^{\gamma} \leqq(\Sigma a)^{\alpha}(\Sigma b)^{\beta}(\Sigma c)^{\gamma}$.
This process also may be repeated, and leads to the general form of H .

We may arrange the induction differently, increasing the number of sets first. The intermediate generalisation ( H for any number of sets of two numbers) is worth separate statement.
40. If $\alpha+\beta+\ldots+\lambda=1$, then

$$
a_{1}^{\alpha} b_{1}^{\beta} \ldots l_{1}^{\lambda}+a_{2}^{\alpha} b_{2}^{\beta} \ldots l_{2}^{\lambda}<\left(a_{1}+a_{2}\right)^{\alpha}\left(b_{1}+b_{2}\right)^{\beta} \ldots\left(l_{1}+l_{2}\right)^{\lambda},
$$

unless $a_{1} / a_{2}=b_{1} / b_{2}=\ldots=l_{1} / l_{2}$ or one of the sets is nul.
(ii) Similarly we can generalise $\mathrm{M}_{0}$ in two directions. On the one hand

$$
\begin{aligned}
& \left\{\left(a_{1}+b_{1}+c_{1}\right)^{r}+\left(a_{2}+b_{2}+c_{2}\right)^{r}\right\}^{1 / r} \\
& \quad \leqq\left(a_{1}^{r}+a_{2}^{r}\right)^{1 / r}+\left\{\left(b_{1}+c_{1}\right)^{r}+\left(b_{2}+c_{2}\right)^{r}\right\}^{1 / r} \\
& \quad \leqq\left(a_{1}^{r}+a_{2}^{r}\right)^{1 / r}+\left(b_{1}^{r}+b_{2}^{r}\right)^{1 / r}+\left(c_{1}^{r}+c_{2}^{r}\right)^{1 / r}
\end{aligned}
$$

and on the other

$$
\begin{aligned}
& \left\{\left(a_{1}+b_{1}\right)^{r}+\left(a_{2}+b_{2}\right)^{r}+\left(a_{3}+b_{3}\right)^{r}\right\}^{1 / r} \\
& \quad \leqq\left[\left\{\left(a_{1}^{r}+a_{2}^{r}\right)^{1 / r}+\left(b_{1}^{r}+b_{2}^{r}\right)^{1 / r}\right\}^{r}+\left(a_{3}+b_{3}\right)^{r}\right]^{1 / r} \\
& \quad \leqq\left(a_{1}^{r}+a_{2}^{r}+a_{3}^{r}\right)^{1 / r}+\left(b_{1}^{r}+b_{2}^{r}+b_{3}^{r}\right)^{1 / r}
\end{aligned}
$$

Repeating and combining these processes, we arrive at the general case.
2.15. Elementary inequalities connected with Theorem 37. We can write $G_{0}$ in the form
or

$$
\begin{gathered}
a^{\alpha}<\{a \alpha+b(1-\alpha)\} b^{\alpha-1} \\
a^{\alpha}-b^{\alpha}<\alpha b^{\alpha-1}(a-b) \quad(0<\alpha<1)
\end{gathered}
$$

which is one case of a system of inequalities prominent in textbooks of analysis. The complete system is stated in Theorem 41 below. The theorem is so important that it is worth while to give a direct proof from first principles which conforms strictly to the criteria of § 1.7.
41. If $x$ and $y$ are positive and unequal, then
(2.15.1) $\quad r x^{r-1}(x-y)>x^{r}-y^{r}>r y^{r-1}(x-y) \quad(r<0$ or $r>1)$,
(2.15.2) $\quad r x^{r-1}(x-y)<x^{r}-y^{r}<r y^{r-1}(x-y) \quad(0<r<1)$.

There is obviously equality when $r=0, r=1$, or $x=y$. We begin by reducing the theorem to one of its cases.
(i) We may suppose $r$ positive. For let us assume that (2.15.1)
has been proved when $r>1$, and that $r<0, r=-s$, so that $s+1>1$. Then

$$
\begin{aligned}
x^{r}-y^{r} & =x^{-s}-y^{-s}=x^{-s} y^{-s-1}\left(y^{s+1}-x^{s} y\right) \\
& =x^{-s} y^{-s-1}\left\{y^{s+1}-x^{s+1}-x^{s}(y-x)\right\}>x^{-s} y^{-s-1} s x^{s}(y-x) \\
& =r y^{r-1}(x-y)
\end{aligned}
$$

The other inequality in (2.15.1) can be treated similarly.
(ii) Let us denote the left- and right-hand inequalities in (2.15.1) by ( $1 a$ ) and ( $1 b$ ) respectively, and similarly for (2.15.2). If we interchange $x$ and $y,(1 b)$ and (2b) become (1 $a$ ) and ( $2 a$ ). It is therefore sufficient to prove ( $1 b$ ) and (2b).
(iii) We may now suppose, on grounds of homogeneity, that $y=1$.

The proof of Theorem 41 is now reduced to that of the next theorem.
42. If $x$ is positive and not equal to 1 , then

$$
\begin{array}{ll}
x^{r}-1>r(x-1) & (r>1)  \tag{2.15.3}\\
x^{r}-1<r(x-1) & (0<r<1)
\end{array}
$$

If in (2.15.3) we write $r=1 / s$ and $x=y^{1 / r}=y^{s}$, it becomes (2.15.4) with $y, s$ for $x, r$. It is therefore sufficient to prove (2.15.3).

If $q$ is an integer greater than $1,{ }^{\text {a }}$ and $y>1$, then

$$
q y^{q}>1+y+\ldots+y^{q-1}=\frac{y^{q}-1}{y-1}>q
$$

If $0<y<1$, the inequalities are reversed. Replacing $y^{q}$ by $x$, we obtain in either case

$$
\begin{equation*}
\frac{x-1}{x}<q\left(x^{1 / q}-1\right)<x-1 \tag{2.15.5}
\end{equation*}
$$

Next, we have

$$
\begin{aligned}
\frac{y^{q+1}-1}{q+1} & -\frac{y^{q}-1}{q}=\frac{y-1}{q(q+1)}\left(q y^{q}-y^{q-1}-y^{q-2}-\ldots-1\right) \\
& =\frac{(y-1)^{2}}{q(q+1)}\left\{y^{q-1}+\left(y^{q-1}+y^{q-2}\right)+\ldots+\left(y^{q-1}+y^{q-2}+\ldots+1\right)\right\}
\end{aligned}
$$

a We abandon here our usual convention concerning the meanings of $q$ and $p$.

The curly bracket contains $\frac{1}{2} q(q+1)$ terms, all of which lie between $y^{q}$ and 1 , so that

$$
\begin{equation*}
\frac{1}{2}(y-1)^{2} \lesseqgtr \frac{y^{q+1}-1}{q+1}-\frac{y^{q}-1}{q} \lesseqgtr \frac{1}{2} y^{q}(y-1)^{2} \quad(y \gtrless 1) ; \tag{2.15.6}
\end{equation*}
$$

and so, if $p$ is any integer greater than $q$,
(2.15.7)

$$
\frac{1}{2}(p-q)(y-1)^{2} \lesseqgtr \frac{y^{p}-1}{p}-\frac{y^{q}-1}{q}>\frac{1}{2}(p-q) y^{p}(y-1)^{2} \quad(y \gtrless 1) .
$$

Now it follows from (2.15.5) that

$$
\frac{(x-1)^{2}}{x^{2}}<q^{2}\left(x^{1 / q}-1\right)^{2}<(x-1)^{2}
$$

if $x>1$, while if $0<x<1$ the inequalities are reversed. Hence, replacing $y^{q}$ by $x$ in (2.15.7), we obtain
(2.15.8)

$$
\frac{p-q}{2 q} \frac{(x-1)^{2}}{x^{2}} \lesseqgtr \frac{x^{p / q}-1}{p / q}-(x-1) \lesseqgtr \frac{p-q}{2 q} x^{p / q}(x-1)^{2} \quad(x \geqslant 1) .
$$

Suppose now that $r>1$. If $r$ is rational, we write $r$ for $p / q$; if $r$ is irrational, we make $p / q \rightarrow r$. In either case we have
$\frac{1}{2}(r-1) \frac{(x-1)^{2}}{x^{2}} \geqq \frac{x^{r}-1}{r}-(x-1) \geqq \frac{1}{2}(r-1) x^{r}(x-1)^{2} \quad(r>1, x \gtrless 1)$, which plainly includes (2.15.3).

This proves Theorems 42 and 41, but it will be useful to have the inequalities corresponding to (2.15.9) when $r<1$. We now replace $y^{p}$ by $x$ in (2.15.7), and use (2.15.5) with $q$ replaced by $p$. We thus obtain
(2.15.10)

$$
\frac{p-q}{2 p} \frac{(x-1)^{2}}{x^{2}} \lesseqgtr x-1-\frac{x^{q / p}-1}{q / p} \lesseqgtr \frac{p-q}{2 p} x(x-1)^{2} \quad(x \gtrless 1),
$$

(2.15.11)

$$
\frac{1}{2}(1-r) \frac{(x-1)^{2}}{x^{2}} \leq x-1-\frac{x^{r}-1}{r} \underset{\underline{\leq}}{\underline{1}} \frac{2}{2}(1-r) x(x-1)^{2} \quad(0<r<1, x \geqslant 1) .
$$

We have made the proof of (2.15.3) rather more elaborate than is necessary, in order to obtain the 'second order' inequalities (2.15.6)(2.15.11), which are interesting in themselves. If we are concerned only to prove (2.15.3), we can argue as follows. Instead of (2.15.6) we write simply
whence

$$
\begin{aligned}
& \frac{y^{q+1}-1}{q+1}>\frac{y^{q}-1}{q} \\
& \frac{y^{p}-1}{p}>\frac{y^{q}-1}{q}
\end{aligned}
$$

if $p$ and $q$ are integers and $p>q$. Hence we obtain (2.15.3) for rational $r$, and so, by a passage to the limit

$$
x^{r}-1 \geqq r(x-1),
$$

for any $r>1$. If now $r$ is irrational, we may write $r=\alpha s$, where $\alpha$ and $s$ are both greater than 1 and $\alpha$ is rational. Then

$$
x^{r}-1=\left(x^{s}\right)^{a}-1>\alpha\left(x^{s}-1\right) \geqq \alpha s(x-1)=r(x-1),
$$

so that (2.15.3) is true generally.
For other proofs of Theorem 41 which satisfy the requirements, see Stolz and Gmeiner (1, 202-208) and Pringsheim (1). Pringsheim uses the result to obtain an elementary proof of $H$. Radon (1,1351) deduces H and M from Theorem 41, but proves this by differential calculus. The proofs of Theorem 41 given in textbooks are usually limited to rational $r$; see for example Chrystal (1, 42-45), Hardy (1, 138).
2.16. Elementary proof of Theorem 3. We have proved incidentally in the last section a number of inequalities sharper than those stated in Theorems 41 and 42. We lay no stress on these, since it is easy to find still more precise inequalities by the aid of the differential calculus (see § 4.2); but it may be interesting to show shortly how they enable us, if we desire, to 'elementarise' the proof of Theorem 3.

We observe first that

$$
\begin{equation*}
a^{r}=1+O(r) \tag{2.16.1}
\end{equation*}
$$

for fixed positive $a$ and small (positive or negative) $r$;

$$
\begin{equation*}
(1+u)^{q}=1+q u+O\left(u^{2}\right) \tag{2.16.2}
\end{equation*}
$$

for fixed $q$ and small $u$; and

$$
\begin{equation*}
\left\{1+O\left(r^{2}\right)\right\}^{1 / r}=1+O(r) \tag{2.16.3}
\end{equation*}
$$

for small $r$. We leave the deduction of these formulae from those of the last section to the reader.

Supposing now that $r$ is small, we have $a_{\nu}^{r}=1+u_{\nu}$, where $u_{\nu}=O(r)$, by (2.16.1), and

$$
a_{\nu}{ }^{q_{\nu}^{r}}=\left(1+u_{\nu}\right)^{q_{\nu}}=1+q_{\nu} u_{\nu}+O\left(r^{2}\right),
$$

by (2.16.2). Hence

$$
\begin{gathered}
\frac{\mathfrak{G}}{\mathfrak{M}_{r}}=\left(\frac{a_{1} q_{1} r}{q_{1} a_{1}^{r}+\ldots a_{n}^{q_{n}^{r}}}\right)^{1 / r}=\left\{\frac{\left(1+u_{n}\right)^{q_{1}} \ldots\left(1+u_{n}\right)^{q_{n}}}{q_{1}\left(1+u_{1}\right)+\ldots+q_{n}\left(1+u_{n}\right)}\right\}^{1 / r} \\
=\left\{\frac{1+q_{1} u_{1}+\ldots+q_{n} u_{n}+O\left(r^{2}\right)}{1+q_{1} u_{1}+\ldots+q_{n} u_{n}}\right\}^{1 / r} \\
=\left\{1+O\left(r^{2}\right)\right\}^{1 / r}=1+O(r) \rightarrow 1 .
\end{gathered}
$$

2.17. Tchebychef's inequality. We know (Theorem 24) that $\mathfrak{M}_{r}(a+b)$ is comparable (§1.6) with $\mathfrak{M}_{r}(a)+\mathfrak{M}_{r}(b)$. It is natural to ask whether $\mathfrak{M r}_{r}(a b)$ is comparable with $\mathfrak{M}_{r}(a) \mathfrak{M}_{r}(b)$. Theorem 43 below shows that this is not so.

We say that ( $a$ ) and (b) are similarly ordered if

$$
\left(a_{\mu}-a_{\nu}\right)\left(b_{\mu}-b_{\nu}\right) \geqq 0,
$$

for all $\mu, \nu$, and oppositely ordered if the inequality is always reversed. It is evident that $(a)$ and $(b)$ are similarly ordered if there is a permutation $\nu_{1}, \nu_{2}, \ldots, \nu_{n}$ of the suffixes such that $a_{\nu_{1}}, a_{\nu_{2}}, \ldots, a_{\nu_{n}}$ and $b_{\nu_{1}}, b_{\nu_{2}}, \ldots, b_{\nu_{n}}$ are both non-decreasing sequences, and oppositely ordered if $a_{\nu_{1}}, \ldots$ is non-increasing and $b_{\nu_{1}}, \ldots$ non-decreasing; and that the converses of these propositions are also true.
43. ${ }^{\text {a }}$ If $r>0$, and ( $a$ ) and (b) are similarly ordered, then

$$
\begin{equation*}
\mathfrak{M}_{r}(a) \mathfrak{M}_{r}(b)<\mathfrak{M}_{r}(a b), \tag{2.17.1}
\end{equation*}
$$

unless all the $a$ or all the $b$ are equal. The inequality is reversed when the sets are oppositely ordered.

It is enough, after (2.2.7), to consider the case $r=1$. Then

$$
\begin{aligned}
\Sigma p \Sigma & \Sigma a b-\Sigma p a \Sigma p b=\Sigma p_{\mu} \Sigma p_{\nu} a_{\nu} b_{\nu}-\Sigma p_{\mu} a_{\mu} \Sigma p_{\nu} b_{\nu} \\
& =\Sigma \Sigma\left(p_{\mu} p_{\nu} a_{\nu} b_{\nu}-p_{\mu} p_{\nu} a_{\mu} b_{\nu}\right)=\Sigma \Sigma\left(p_{\nu} p_{\mu} a_{\mu} b_{\mu}-p_{\nu} p_{\mu} a_{\nu} b_{\mu}\right) \\
\quad & =\frac{1}{2} \Sigma \Sigma\left(p_{\mu} p_{\nu} a_{\nu} b_{\nu}-p_{\mu} p_{\nu} a_{\mu} b_{\nu}+p_{\nu} p_{\mu} a_{\mu} b_{\mu}-p_{\nu} p_{\mu} a_{\nu} b_{\mu}\right) \\
& =\frac{1}{2} \Sigma \Sigma p_{\mu} p_{\nu}\left(a_{\mu}-a_{\nu}\right)\left(b_{\mu}-b_{\nu}\right) \geqq 0
\end{aligned}
$$

or

$$
\mathfrak{A}(a) \mathfrak{A}(b) \leqq \mathfrak{A}(a b)
$$

if the series are similarly ordered.
We can determine the cases of equality as follows. Suppose, as
a The integral analogue is due to Tchebychef. See Hermite (1, 46-47), Franklin (1), Jensen (1), and Theorem 236. When $r=1, \mathfrak{M}_{r}=\mathfrak{A}$, the inequality holds for any real and similarly ordered $a, b$.
we may in virtue of the remarks earlier in the section, that ( $a$ ) and (b) are non-decreasing. The double sum contains a term

$$
p_{1} p_{n}\left(a_{1}-a_{n}\right)\left(b_{1}-b_{n}\right)
$$

and this can vanish only if all $a$ or all $b$ are equal.
An immediate deduction is that

$$
\mathfrak{M}_{r}(a) \mathfrak{M r}_{r}(b) \ldots \mathfrak{M r}_{r}(l)<\mathfrak{M r}_{r}(a b \ldots l)
$$

if $r$ is positive and $(a),(b), \ldots,(l)$ are all similarly ordered; in particular

$$
\mathfrak{M}_{r}(a)<\mathfrak{M r}_{m r}(a)
$$

if $m$ is an integer greater than 1. This includes Theorem 6 and is included in Theorem 16.

The question asked at the beginning of this section is included in the more general question settled by the next theorem.
44. A necessary and sufficient condition that $\mathfrak{M}_{r}(a b \ldots l)$ and $\mathfrak{M}_{s}(a) \mathfrak{M l}_{t}(b) \ldots \mathfrak{M}_{v}(l)$, where $r, s, \ldots, v$ are positive, should be comparable, is that

$$
\begin{equation*}
\frac{1}{r} \geqq \frac{1}{s}+\frac{1}{t}+\ldots+\frac{1}{v} \tag{2.17.2}
\end{equation*}
$$

in which case

$$
\begin{equation*}
\mathfrak{M d}_{r}(a b \ldots l) \leqq \mathfrak{M}_{s}(a) \mathfrak{M}_{t}(b) \ldots \mathfrak{M}_{v}(l) . \tag{2.17.3}
\end{equation*}
$$

The sufficiency of the condition follows at once from Theorems 12 and 16. If we take every set in (2.17.3) to be ( $1,0,0, \ldots, 0$ ), we see at once that (2.17.2) must be satisfied. A general inequality opposite to (2.17.3) is impossible for any $r, s, \ldots$, since $a_{\nu} b_{\nu} \ldots l_{\nu}$ may vanish for every $\nu$ and yet the right-hand side be positive.
2.18. Muirhead's theorem. In this and the four succeeding sections we suppose the $a$ to be strictly positive. We denote by

$$
\Sigma!F\left(a_{1}, a_{2}, \ldots, a_{n}\right)
$$

the sum of the $n!$ terms obtained from $F\left(a_{1}, a_{2}, \ldots, a_{n}\right)$ by the possible permutations of the $a$. We shall be concerned only with the special case

$$
F\left(a_{1}, a_{2}, \ldots, a_{n}\right)=a_{1}{ }^{\alpha_{1}} a_{2}{ }^{\alpha_{2}} \ldots a_{n}{ }^{\alpha_{n}} \quad\left(a_{\nu}>0, \alpha_{\nu} \geqq 0\right) .
$$

We write

$$
[\alpha]=\left[\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}\right]=\frac{1}{n!} \Sigma!a_{1}^{\alpha_{1}} a_{2}^{\alpha_{2}} \ldots a_{n}^{\alpha_{n}}
$$

It is plain that $[\alpha]$ is unaltered by any permutation of the $\alpha$, so that we may regard two sets of $\alpha$ as the same if they differ only in arrangement. We may describe a mean value of the type $[\alpha]$ as a symmetrical mean.

In particular

$$
\begin{aligned}
& {[1,0,0, \ldots, 0]=\frac{(n-1)!}{n!}\left(a_{1}+a_{2}+\ldots+a_{n}\right)=\mathfrak{A}(a),} \\
& {\left[\frac{1}{n}, \frac{1}{n}, \ldots, \frac{1}{n}\right]=\frac{n!}{n!} a_{1}^{1 / n} a_{2}^{1 / n} \ldots a_{n}^{1 / n}=\mathfrak{S}(a)}
\end{aligned}
$$

the arithmetic and geometric means with unit weights. When $\alpha_{1}+\alpha_{2}+\ldots+\alpha_{n}=1,[\alpha]$ is a common generalisation of $\mathfrak{A}(a)$ and (3) (a).

In general $\left[\alpha^{\prime}\right]$ is not comparable with $[\alpha]$ in the sense of $\S 1.6$. The problem solved in this and the next two sections is that of determining conditions for comparability.

We say that ( $\alpha^{\prime}$ ) is majorised by ( $\alpha$ ), and write

$$
\left(\alpha^{\prime}\right) \prec(\alpha),
$$

when the ( $\alpha$ ) and ( $\alpha^{\prime}$ ) can be arranged so as to satisfy the following three conditions:

$$
\begin{align*}
& \alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\ldots+\alpha_{n}^{\prime}=\alpha_{1}+\alpha_{2}+\ldots+\alpha_{n}  \tag{2.18.1}\\
& \alpha_{1}^{\prime} \geqq \alpha_{2}^{\prime} \geqq \ldots \geqq \alpha_{n}^{\prime}, \quad \alpha_{1} \geqq \alpha_{2} \geqq \ldots \geqq \alpha_{n}  \tag{2.18.2}\\
& \alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\ldots+\alpha_{\nu}^{\prime} \leqq \alpha_{1}+\alpha_{2}+\ldots+\alpha_{\nu} \quad(1 \leqq v<n) . \tag{2.18.3}
\end{align*}
$$

The second condition is in itself no restriction, since we may rearrange ( $\alpha^{\prime}$ ) and ( $\alpha$ ) in any order, but it is essential to the statement of the third. It is plain that $(\alpha)<(\alpha)$.
45. A necessary and sufficient condition that $\left[\alpha^{\prime}\right]$ should be comparable with $[\alpha]$, for all positive values of the $a$, is that one of $\left(\alpha^{\prime}\right)$ and $(\alpha)$ should be majorised by the other. If $\left(\alpha^{\prime}\right)<(\alpha)$ then

$$
\begin{equation*}
\left[\alpha^{\prime}\right] \leqq[\alpha] . \tag{2.18.4}
\end{equation*}
$$

There is equality only when ( $\alpha^{\prime}$ ) and ( $\alpha$ ) are identical or when all the a are equala.

[^11]2.19. Proof of Muirhead's theorem. (1) The condition is necessary. Suppose, as we may, that (2.18.2) is satisfied, and that (2.18.4) holds for all positive $a$. Taking all the $a$ equal to $x$, we obtain
$$
x^{\Sigma \alpha^{\prime}}=\left[\alpha^{\prime}\right] \leqq[\alpha]=x^{\Sigma \alpha} .
$$

This can only be true both for large and for small $x$ if (2.18.1) is true.

Next, take

$$
a_{1}=a_{2}=\ldots=a_{\nu}=x, \quad a_{\nu+1}=\ldots=a_{n}=1
$$

$x$ being large. Since $\left(\alpha^{\prime}\right)$ and $(\alpha)$ are in descending order, the indices of the highest powers of $x$ in $\left[\alpha^{\prime}\right]$ and $[\alpha]$ are

$$
\alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\ldots+\alpha_{\nu}^{\prime}, \quad \alpha_{1}+\alpha_{2}+\ldots+\alpha_{\nu}
$$

respectively. It is plain that the first cannot exceed the second, and this proves (2.18.3).
(2) The condition is sufficient. The proof of this is rather more troublesome, and we require a new definition and two lemmas.

We define a special type of linear transformation of the $\alpha$, which we call a transformation $T$, as follows. Suppose that $\alpha_{k}$ and $\alpha_{l}$ are two unequal $\alpha$, the first being the greater; we may write

$$
\begin{equation*}
\alpha_{k}=\rho+\tau, \quad \alpha_{l}=\rho-\tau \quad(0<\tau \leqq \rho) \tag{2.19.1}
\end{equation*}
$$

If now

$$
\begin{equation*}
0 \leqq \sigma<\tau \leqq \rho \tag{2.19.2}
\end{equation*}
$$

then a transformation $T$ is defined by

$$
\left\{\begin{array}{l}
\alpha_{k}^{\prime}=\rho+\sigma=\frac{\tau+\sigma}{2 \tau} \alpha_{k}+\frac{\tau-\sigma}{2 \tau} \alpha_{l}  \tag{2.19.3}\\
\alpha_{l}^{\prime}=\rho-\sigma=\frac{\tau-\sigma}{2 \tau} \alpha_{k}+\frac{\tau+\sigma}{2 \tau} \alpha_{l} \\
\alpha_{\nu}^{\prime}=\alpha_{\nu} \quad(\nu \neq k, \nu \neq l)
\end{array}\right.
$$

If ( $\alpha^{\prime}$ ) arises from ( $\alpha$ ) by a transformation $T$, we write $\alpha^{\prime}=T \alpha$. The definition does not necessarily imply that either the $\alpha$ or the $\alpha^{\prime}$ are in decreasing order.

It is plain that the sufficiency of our condition for comparability will be established, and that we shall also have proved what
is stated in Theorem 45 about the case of equality, if we have proved the two following lemmas.
Lemma 1. If $\alpha^{\prime}=T \alpha$ then $\left[\alpha^{\prime}\right] \leqq[\alpha]$, with equality only when all a are equal.

Lemma 2. If $\left(\alpha^{\prime}\right) \prec(\alpha)$, but ( $\alpha^{\prime}$ ) is not identical with ( $\alpha$ ), then $\left.{ }^{( } \alpha^{\prime}\right)$ can be derived from ( $\alpha$ ) by the successive application of a finite number of transformations $T$.

Proof of Lemma 1. We may rearrange ( $\alpha$ ) and ( $\alpha^{\prime}$ ) so that $k=1, l=2$. Then

$$
\begin{aligned}
& \text { (2.19.4) } n!2[\alpha]-n!2\left[\alpha^{\prime}\right] \\
& =n!2\left[\rho+\tau, \rho-\tau, \alpha_{3}, \ldots\right]-n!2\left[\rho+\sigma, \rho-\sigma, \alpha_{3}, \ldots\right] \\
& =\Sigma!a_{3}{ }^{\alpha_{3}} \ldots a_{n}{ }^{\alpha_{n}}\left(a_{1}{ }^{\rho+\tau} a_{2}{ }^{\rho-\tau}+a_{1}{ }^{\rho-\tau} a_{2}{ }^{\rho+\tau}-a_{1}{ }^{\rho+\sigma} a_{2}{ }^{\rho-\sigma}-a_{1}{ }^{\rho-\sigma} a_{2}{ }^{\rho+\sigma}\right) \\
& =\Sigma!\left(a_{1} a_{2}\right)^{\rho-\tau} a_{3}^{\alpha_{3}} \ldots a_{n}^{\alpha_{n}}\left(a_{1}{ }^{\tau+\sigma}-a_{2}{ }^{\tau+\sigma}\right)\left(a_{1}{ }^{\tau-\sigma}-a_{2}{ }^{\tau-\sigma}\right) \geqq 0 \text {, }
\end{aligned}
$$

with equality only when all the $a$ are equal.
Proof of Lemma 2. We suppose that the condition (2.18.2) is satisfied, and call the number of the differences $\alpha_{\nu}-\alpha_{\nu}{ }^{\prime}$ which are not zero the discrepancy of ( $\alpha$ ) and ( $\alpha^{\prime}$ ); if the discrepancy is zero the sets are identical. We prove the lemma by induction, assuming it to be true when the discrepancy is less than $r$ and proving that it is then true when the discrepancy is $r$.

Suppose then that $\left(\alpha^{\prime}\right) \prec(\alpha)$ and that the discrepancy is $r>0$. Since, by (2.18.1), $\Sigma\left(\alpha_{\nu}-\alpha_{\nu}{ }^{\prime}\right)=0$, and not all of these differences are zero, there must be positive and negative differences; and, by (2.18.3) the first which is not zero must be positive. We can therefore find $k$ and $l$ so that

$$
\begin{equation*}
\alpha_{k}^{\prime}<\alpha_{k}, \quad \alpha_{k+1}^{\prime}=\alpha_{k+1}, \ldots, \quad \alpha_{l-1}^{\prime}=\alpha_{l-1}, \quad \alpha_{l}^{\prime}>\alpha_{l} . a^{a} \tag{2.19.5}
\end{equation*}
$$

We take $\alpha_{k}=\rho+\tau, \alpha_{l}=\rho-\tau$, as in (2.19.1), and define $\sigma$ by

$$
\begin{equation*}
\sigma=\operatorname{Max}\left(\left|\alpha_{k}^{\prime}-\rho\right|,\left|\alpha_{l}^{\prime}-\rho\right|\right) . \tag{2.19.6}
\end{equation*}
$$

Then $0<\tau \leqq \rho$, since $\alpha_{k}>\alpha_{l}$. Also one or other ${ }^{b}$ of

$$
\alpha_{l}^{\prime}-\rho=-\sigma, \quad \alpha_{k}^{\prime}-\rho=\sigma
$$

[^12]is true, since $\alpha_{k}^{\prime} \geqq \alpha_{l}^{\prime}$; and $\sigma<\tau$, since $\alpha_{k}^{\prime}<\alpha_{k}$ and $\alpha_{l}^{\prime}>\alpha_{l}$. Hence
$$
0 \leqq \sigma<\tau \leqq \rho,
$$
as in (2.19.2).
We now write
\[

$$
\begin{equation*}
\alpha_{k}^{\prime \prime}=\rho+\sigma, \quad \alpha_{l}^{\prime \prime}=\rho-\sigma, \quad \alpha_{\nu}^{\prime \prime}=\alpha_{\nu} \quad(\nu \neq k, \nu \neq l) \tag{2.19.7}
\end{equation*}
$$

\]

If $\alpha_{k}{ }^{\prime}-\rho=\sigma, \alpha_{k}{ }^{\prime \prime}=\alpha_{k}{ }^{\prime}$; if $\alpha_{l}{ }^{\prime}-\rho=-\sigma, \alpha_{l}{ }^{\prime \prime}=\alpha_{l}{ }^{\prime}$. ${ }^{\text {a }}$ Since the pairs $\alpha_{k}, \alpha_{k}^{\prime}$ and $\alpha_{l}, \alpha_{l}^{\prime}$ each contribute a unit to the discrepancy $r$ between ( $\alpha^{\prime}$ ) and ( $\alpha$ ), the discrepancy between ( $\alpha^{\prime}$ ) and ( $\alpha^{\prime \prime}$ ) is smaller, being $r-1$ or $r-2$.

Next, comparing (2.19.7) with (2.19.3), and observing that (2.19.2) is satisfied, we see that ( $\alpha^{\prime \prime}$ ) arises from ( $\alpha$ ) by a transformation $T$.

Finally, $\left(\alpha^{\prime}\right)$ is majorised by ( $\alpha^{\prime \prime}$ ). To prove this we must verify that the conditions corresponding to (2.18.1), (2.18.2) and (2.18.3), with $\alpha^{\prime \prime}$ for $\alpha$, are satisfied. For the first, we have

$$
\begin{equation*}
\alpha_{l c}^{\prime \prime}+\alpha_{l}^{\prime \prime}=2 \rho=\alpha_{k}+\alpha_{l}, \quad \Sigma \alpha^{\prime}=\Sigma \alpha=\Sigma \alpha^{\prime \prime} \tag{2.19.8}
\end{equation*}
$$

For the second, we observe first that

$$
\begin{aligned}
\alpha_{k}^{\prime} \leqq \rho+\left|\alpha_{k}^{\prime}-\rho\right| \leqq \rho+\sigma=\alpha_{k}^{\prime \prime}, \\
\alpha_{l}^{\prime} \leqq \rho-\left|\alpha_{l}^{\prime}-\rho\right| \geqq \rho-\sigma=\alpha_{l}^{\prime \prime}
\end{aligned}
$$

and so, by (2.19.5),

$$
\begin{aligned}
& \alpha^{\prime \prime}{ }_{k-1}=\alpha_{k-1} \geqq \alpha_{k}=\rho+\tau>\rho+\sigma=\alpha^{\prime \prime}{ }_{k} \geqq \alpha_{k}^{\prime} \geqq \alpha^{\prime}{ }_{k+1}=\alpha_{k+1}=\alpha^{\prime \prime}{ }_{k+1}, \\
& \alpha^{\prime \prime}{ }_{l-1}=\alpha_{l-1}=\alpha_{l-1}^{\prime} \geqq \alpha^{\prime} \geqq \alpha_{l}^{\prime \prime}=\rho-\sigma>\rho-\tau=\alpha_{l} \geqq \alpha_{l+1}=\alpha^{\prime \prime}{ }_{l+1}
\end{aligned}
$$

and the inequalities affecting the $\alpha^{\prime \prime}$ are those required. Finally, we have to prove that

$$
\alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\ldots+\alpha_{\nu}^{\prime} \leqq \alpha_{1}^{\prime \prime}+\alpha_{2}^{\prime \prime}+\ldots+\alpha_{\nu}^{\prime \prime}
$$

Now this is true if $\nu<k$ or $\nu \geqq l$, by (2.19.7) and (2.18.3); it is true for $\nu=k$, because it is true for $\nu=k-1$ and $\alpha_{k}{ }^{\prime} \leqq \alpha_{k}{ }^{\prime \prime}$; and it is true for $k<\nu<l$ because it is true for $\nu=k$ and the intervening $\alpha^{\prime}$ and $\alpha^{\prime \prime}$ are identical.

We have thus proved that ( $\alpha^{\prime}$ ) is majorised by ( $\alpha^{\prime \prime}$ ), a set arising from ( $\alpha$ ) by a transformation $T$ and having a discrepancy from $\left(\alpha^{\prime}\right)$ less than $r$. This proves Lemma 2 and so completes the proof of Theorem 45.b
a Again, both these equations may be true.
b For another proof, see Theorems 74 and 75.
2.20. An alternative theorem. We shall say that ( $\alpha^{\prime}$ ) is an average of $(\alpha)$ if there are $n^{2}$ numbers $p_{\mu \nu}$ such that

$$
\begin{equation*}
p_{\mu \nu} \geqq 0, \quad \sum_{\mu=1}^{n} p_{\mu \nu}=1, \quad \sum_{\nu=1}^{n} p_{\mu \nu}=1 \tag{2.20.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\alpha_{\mu}^{\prime}=p_{\mu 1} \alpha_{1}+p_{\mu 2} \alpha_{2}+\ldots+p_{\mu n} \alpha_{n} . \tag{2.20.2}
\end{equation*}
$$

Since the conditions (2.20.1) are not affected by a permutation of the $\mu$ or the $\nu$, the definition is, like that of $\S 2.18$, independent of the order of the $\alpha$ or the $\alpha^{\prime}$. The equations (2.19.3) show that $\left(\rho+\sigma, \rho-\sigma, \alpha_{3}, \ldots\right)$ is an average of $\left(\rho+\tau, \rho-\tau, \alpha_{3}, \ldots\right)$ when (2.19.2) is satisfied.

The last two conditions (2.20.1) may also be stated as follows: $\Sigma \alpha^{\prime}$, when expressed as a function of the $\alpha$, is identical with $\Sigma \alpha$, and every $\alpha^{\prime}$ is 1 if every $\alpha$ is 1 . From this it follows that the relationship is transitive; if ( $\alpha^{\prime}$ ) is an average of ( $\alpha$ ), and ( $\alpha^{\prime \prime}$ ) of $\left(\alpha^{\prime}\right)$, then ( $\alpha^{\prime \prime}$ ) is an average of ( $\alpha$ ). And from this and Lemma 2 of $\S 2.19$ it follows that if $\left(\alpha^{\prime}\right) \prec(\alpha)$ then $\left(\alpha^{\prime}\right)$ is an average of $(\alpha)$.

The converse is also true. For suppose that (2.20.1) and (2.20.2) are satisfied. Then (2.18.1) follows by addition of the equations (2.20.2). Finally, if we suppose ( $\alpha$ ) and ( $\alpha^{\prime}$ ) in descending order, and write

$$
p_{1 \nu}+p_{2 \nu}+\ldots+p_{m \nu}=k_{\nu},
$$

we have $k_{\nu} \leqq 1$ and $\Sigma k_{\nu}=m$, by (2.20.1); and so

$$
\begin{gathered}
\alpha_{1}^{\prime}+\alpha_{2}^{\prime}+\ldots+\alpha_{m}{ }^{\prime} \leqq k_{1} \alpha_{1}+\ldots+k_{m-1} \alpha_{m-1}+\left(m-k_{1}-\ldots-k_{m-1}\right) \alpha_{m} \\
\leqq\left(\alpha_{1}-\alpha_{m}\right)+\ldots+\left(\alpha_{m-1}-\alpha_{m}\right)+m \alpha_{m}=\alpha_{1}+\alpha_{2}+\ldots+\alpha_{m},
\end{gathered}
$$

which is (2.18.3).
We have therefore proved the two following theorems.
46. A necessary and sufficient condition that ( $\alpha^{\prime}$ ) should be an average of $(\alpha)$ is that $\left(\alpha^{\prime}\right) \prec(\alpha)$.
47. A necessary and sufficient condition that $\left[\alpha^{\prime}\right]$ should be comparable with $[\alpha]$ is that one of $\left(\alpha^{\prime}\right)$ and ( $\alpha$ ) should be an average of the other. If $\left(\alpha^{\prime}\right)$ is an average of $(\alpha)$ then $\left[\alpha^{\prime}\right] \leqq[\alpha]$, with equality only as in Theorem 45.
2.21. Further theorems on symmetrical means. (1) Theorems 45 and 47 fulfil two purposes. First, either theorem gives a
simple criterion for deciding whether two means $[\alpha]$ and $\left[\alpha^{\prime}\right]$ are or are not comparable. Secondly, the proof of Theorem 45 shows us how, by repeated application of the transformation (2.19.3) and the formula (2.19.4), to decompose the difference of two comparable means into a sum of obviously positive terms. We obtain, for example, a new and interesting proof of the theorem of the arithmetic and geometric means (with unit weights); in fact

$$
\begin{aligned}
\mathfrak{A}\left(a^{n}\right)- & \mathscr{S}\left(a^{n}\right)=[n, 0,0, \ldots, 0]-[1,1, \ldots, 1] \\
& \quad([n, 0,0, \ldots, 0]-[n-1,1,0, \ldots, 0]) \\
& +([n-1,1,0, \ldots, 0]-[n-2,1,1,0, \ldots, 0]) \\
& +([n-2,1,1,0, \ldots, 0]-[n-3,1,1,1,0, \ldots, 0])+\ldots \\
=\frac{1}{2(n!)} & \left\{\Sigma!\left(a_{1}{ }^{n-1}-a_{2}{ }^{n-1}\right)\left(a_{1}-a_{2}\right)+\Sigma!\left(a_{1}{ }^{n-2}-a_{2}{ }^{n-2}\right)\left(a_{1}-a_{2}\right) a_{3}\right. \\
& \left.+\Sigma!\left(a_{1}{ }^{n-3}-a_{2}{ }^{n-3}\right)\left(a_{1}-a_{2}\right) a_{3} a_{4}+\ldots\right\} .
\end{aligned}
$$

Since

$$
\left(a_{r}-a_{s}\right)\left(a_{r}^{\nu}-a_{s}^{\nu}\right)>0
$$

unless $a_{r}=a_{s}$, the theorem follows ${ }^{\text {a }}$.
(2) 48. If $\alpha_{1}+\alpha_{2}+\ldots+\alpha_{n}=1$, then

$$
\mathfrak{G}(a)<[\alpha]<\mathfrak{A}(a),
$$

unless $[\alpha]$ is $\mathfrak{G}(a)$ or $\mathfrak{A}(a)$, or all the $a$ are equal.
This theorem ${ }^{\text {b }}$ shows that all the $[\alpha]$ of degree of homogeneity 1 are comparable with $\mathfrak{G}(a)$ and $\mathfrak{A}(a)$, though not in general comparable among themselves. To prove it we apply Theorem 47; since

$$
\frac{1}{n}=\frac{\alpha_{1}}{n}+\frac{\alpha_{2}}{n}+\ldots+\frac{\alpha_{n}}{n}
$$

and

$$
\alpha_{\mu}=\alpha_{\mu} \cdot 1+\alpha_{\mu+1} \cdot 0+\ldots+\alpha_{n} \cdot 0+\alpha_{1} \cdot 0+\ldots+\alpha_{\mu-1} \cdot 0,
$$ $(1 / n, 1 / n, \ldots, 1 / n)$ is an average of $(\alpha)$ and ( $\alpha$ ) an average of $(1,0, \ldots, 0)$. Or we may deduce Theorem 48 directly from Theorem 45.

(3) We add two further theorems of a similar character, with indications only of the proofs.
49. If $0<\sigma \leqq 1$, then a necessary and sufficient condition that $\left[\alpha^{\prime}\right] \leqq[\alpha]^{\sigma}$ is that $\left(\alpha^{\prime}\right)<(\sigma \alpha)$. If $\sigma>1$, the condition is necessary but not sufficient.

[^13][To prove the condition necessary, follow the line of § 2.19 (1). To prove it sufficient, combine Theorems 45 and 11. As an example
$$
[r, 0,0, \ldots] \leqq[s, 0,0, \ldots]^{r / s} \quad(0<r<s) ;
$$
this is $\mathfrak{M}_{r}(a) \leqq \mathfrak{M}_{s}(a)$ (Theorem 16), with unit weights. The same example shows that the condition is no longer sufficient when $\sigma>1$.]
50. If $r, p, a$ are positive and
$$
T_{r}=\Sigma p_{\nu} a_{\nu}^{r}=\mathfrak{T}_{r}^{r}
$$
(in the notation of § 2.10 (iv)), then a necessary and sufficient condition that
$$
T_{a_{1}^{\prime}} T_{a_{2}^{\prime}} \ldots T_{a_{n^{\prime}}} \leqq T_{a_{4}} T_{a_{2}} \ldots T_{a_{n}}
$$
for all $a$ and $p$ is that $\left(\alpha^{\prime}\right)<(\alpha)$.
[The necessity of the condition may be established as before. To prove it sufficient we use Theorem 46 and Hölder's inequality, which give
\[

$$
\begin{aligned}
T_{a_{\mu}^{\prime}} & =T_{s_{1} a_{1}+s_{\mu 2} a_{2}+\cdots+s_{\mu m} a_{m}} \\
& \leqq\left(T_{a_{1}}\right)^{s_{11}}\left(T_{a_{2}}\right)^{s_{\mu 2}} \ldots\left(T_{a_{m}}\right)^{s_{\mu m}}:
\end{aligned}
$$
\]

we have changed the notation slightly in order to avoid conflict with that of § 2.10. The result follows by multiplication.]

### 2.22. The elementary symmetric functions of $n$ positive

 numbers. If$$
\begin{aligned}
\left(x+a_{1}\right)\left(x+a_{2}\right) \ldots(x & \left.+a_{n}\right)=x^{n}+c_{1} x^{n-1}+c_{2} x^{n-2}+\ldots+c_{n} \\
& =x^{n}+\binom{n}{1} p_{1} x^{n-1}+\binom{n}{2} p_{2} x^{n-2}+\ldots+p_{n}
\end{aligned}
$$

then $c_{r}$ is the $r$ th elementary symmetric function of the $a$, i.e. the sum of the products, $r$ at a time, of different $a$, and $p_{r}$ the average of these products. In this section we consider two well-known theorems concerning the $p_{r}$. We write $c_{0}=p_{0}=1$.

In the notation of $\S 2.18$

$$
\begin{aligned}
c_{r} & =\frac{1}{r!(n-r)!} \Sigma!a_{1} a_{2} \ldots a_{r} \\
p_{r} & =\frac{r!(n-r)!}{n!} c_{r}=[1,1, \ldots, 1,0,0, \ldots, 0]
\end{aligned}
$$

there being $r$ l's and $n-r 0^{\prime}$ 's. Also $p_{1}=\mathfrak{A}(a)$ and $p_{n}=\mathscr{S N}^{n}(a)$, with unit weights. The different $p_{r}$, being of different degrees, are not comparablea; but they are connected by non-linear inequalities.
a This is a trivial case of Theorem 45.
51. $p_{r-1} p_{r+1}<p_{r}^{2}(1 \leqq r<n)$, unless all the $a$ are equal.
52. $p_{1}>p_{2^{\frac{1}{2}}}>p_{3^{\frac{b}{3}}}>\ldots>p_{n}^{1 / n}$, unless all the a are equal.

Theorem 51, which was stated by Newton ${ }^{\text {a }}$, is actually true for real, not necessarily positive, $a$; and we shall give a proof of the more general theorem, depending on the methods of the differential calculus, in §4.3. Theorem 52 is due to Maclaurin ${ }^{\text {b }}$.

Theorem 52 is a corollary of Theorem 51, since

$$
\left(p_{0} p_{2}\right)\left(p_{1} p_{3}\right)^{2}\left(p_{2} p_{4}\right)^{3} \ldots\left(p_{r-1} p_{r+1}\right)^{r}<p_{1}^{2} p_{2}^{4} p_{3}^{6} \ldots p_{r}^{2 r}
$$

gives $p_{r+1}^{r}<p_{r}^{r+1}$ or

$$
p_{r}^{1 / r}>p_{r+1}^{1 /(r+1)}
$$

This remark, together with the proof of §4.3, disposes of the theorems, but it is interesting to consider proofs of them by the methods of this chapter.
(i) Proof of Theorem 52 by the method of §2.6 (iii). We begin by proving a theorem similar to but weaker than Theorem 51.
53. ${ }^{\mathrm{c}} c_{r-1} c_{r+1}<c_{r}^{2}$.

This theorem is weaker than Theorem 51, since $p_{r-1} p_{r+1}<p_{r}^{2}$ is

$$
\frac{(r+1)(n-r+1)}{r(n-r)} c_{r-1} c_{r+1}<c_{r}^{2} .
$$

To prove it we observe that a typical term in $c_{r-1} c_{r+1}-c_{r}^{2}$ is

$$
a_{1}^{2} a_{2}^{2} \ldots a_{r-s}^{2} a_{r-s+1} \ldots a_{r+s}
$$

and that this occurs with the coefficient

$$
\binom{2 s}{s-1}-\binom{2 s}{s}<0
$$

From Theorem 53 it follows that

$$
\begin{equation*}
c_{r-1} c_{s}<c_{r} c_{s-1} \tag{2.22.1}
\end{equation*}
$$

if $r<s$.
${ }^{\text {a }}$ Newton (1, 173). See also Maclaurin (2).
${ }^{\text {b }}$ Maclaurin (2). See also Schlömilch (1). The inequality $p_{1}>p_{n}{ }^{1 / n}$ is a case of Theorem 9.
c The theorem is stated, like Theorems 51 and 52, for positive $a$. It remains true, as the proof shows, for non-negative $a$, unless $c_{r}=0$ (i.e. unless all but $r-1$ of the $a$ are 0 ).

We can now prove Theorem 52 as follows. If the $a$ are not all equal, let $a_{1}=\operatorname{Min} a, a_{2}=\operatorname{Max} a$. Then

$$
\begin{gather*}
a_{1}<\alpha_{1}<a_{2}  \tag{2.22.2}\\
\alpha_{1}=p_{\mu}^{1 / \mu}
\end{gather*}
$$

where
We replace $a_{1}$ and $a_{2}$ by $\alpha_{1}$ and $\alpha_{2}$, choosing $\alpha_{2}$ so that $p_{\mu}$ shall be unaltered, and prove that any $p_{\nu}$ for which $\nu>\mu$ is increased by the substitution. The result will then follow as in $\S 2.6$ (iii).

We have

$$
\binom{n}{\mu} p_{\mu}=c_{\mu}=a_{1} a_{2} c_{\mu-2}^{\prime}+\left(a_{1}+a_{2}\right) c_{\mu-1}^{\prime}+c_{\mu}^{\prime}
$$

where $c_{r}{ }^{\prime}$ is the $c_{r}$ formed from the $n-2$ numbers other than $a_{1}$ and $a_{2}$. Since $p_{\mu}$ is to be unchanged

$$
a_{1} a_{2} c_{\mu-2}^{\prime}+\left(a_{1}+a_{2}\right) c_{\mu-1}^{\prime}+c_{\mu}^{\prime}=\alpha_{1} \alpha_{2} c_{\mu-2}^{\prime}+\left(\alpha_{1}+\alpha_{2}\right) c_{\mu-1}^{\prime}+c_{\mu}^{\prime}
$$

(2.22.3) $\quad\left(\alpha_{1} \alpha_{2}-a_{1} a_{2}\right) c_{\mu-2}^{\prime}=-\left(\alpha_{1}+\alpha_{2}-a_{1}-a_{2}\right) c_{\mu-1}^{\prime}$,
(2.22.4) $\quad\left(\alpha_{1} c_{\mu-2}^{\prime}+c_{\mu-1}^{\prime}\right) \alpha_{2}=a_{1} a_{2} c_{\mu-2}^{\prime}+\left(a_{1}+a_{2}-\alpha_{1}\right) c_{\mu-1}^{\prime}$.

The value of $\alpha_{2}$ defined by (2.22.4) is positive because of (2.22.2).
Also, if $p_{v}$ becomes $p_{v}^{*}$,

$$
\binom{n}{\nu}\left(p_{\nu}^{*}-p_{\nu}\right)=\left(\alpha_{1} \alpha_{2}-a_{1} a_{2}\right) c_{\nu-2}^{\prime}+\left(\alpha_{1}+\alpha_{2}-a_{1}-a_{2}\right) c_{\nu-1}^{\prime}
$$

and so $p_{\nu}^{*}-p_{\nu}$ has the sign of

$$
\left(\alpha_{1}+\alpha_{2}-a_{1}-a_{2}\right)\left(\frac{c_{\nu-1}^{\prime}}{c_{\nu-2}^{\prime}}-\frac{c_{\mu-1}^{\prime}}{c_{\mu-2}^{\prime}}\right)
$$

The second factor is negative, by (2.22.1); and, by (2.22.3),

$$
\begin{aligned}
\operatorname{sgn}\left(\alpha_{1}+\alpha_{2}-a_{1}\right. & \left.-a_{2}\right)=\operatorname{sgn}\left(a_{1} a_{2}-\alpha_{1} \alpha_{2}\right) \\
& =\operatorname{sgn}\left\{\alpha_{1}\left(\alpha_{1}+\alpha_{2}-a_{1}-a_{2}\right)+a_{1} a_{2}-\alpha_{1} \alpha_{2}\right\} \\
& =\operatorname{sgn}\left\{\left(\alpha_{1}-a_{1}\right)\left(\alpha_{1}-a_{2}\right)\right\}=-1,
\end{aligned}
$$

by (2.22.2). Hence $p_{v}^{*}>p_{\nu}$, which proves the theorem.
(ii) Proof of Theorem 51 by induction ${ }^{\text {a }}$. Suppose that Theorem 51 has been proved for $n-1$ numbers, $a_{1}, a_{2}, \ldots, a_{n-1}$, and that
${ }^{\text {a }}$ This proof was communicated to us independently by Messrs A. L. Dixon, A. E. Jolliffe, and M. H. A. Newman.
$c_{r}^{\prime}, p_{r}^{\prime}$ are the $c_{r}, p_{r}$ formed from these $n-1$ numbers, which we shall at first suppose not all equal. Then
and so

$$
\begin{gathered}
c_{r}=c_{r}^{\prime}+a_{n} c_{r-1}^{\prime} \\
p_{r}=\frac{n-r}{n} p_{r}^{\prime}+\frac{r}{n} a_{n} p_{r-1}^{\prime}
\end{gathered}
$$

Hence we deduce

$$
\begin{aligned}
& \qquad n^{2}\left(p_{r-1} p_{r+1}-p_{r}^{2}\right)=A+B a_{n}+C a_{n}^{2} \\
& \text { where } \begin{array}{r}
A=\left\{(n-r)^{2}-1\right\} p_{r-1}^{\prime} p_{r+1}^{\prime}-(n-r)^{2} p_{r}^{\prime 2} \\
B=(n-r+1)(r+1) p_{r-1}^{\prime} p_{r}^{\prime}+(n-r-1)(r-1) p_{r-2}^{\prime} p_{r+1}^{\prime} \\
\\
-2 r(n-r) p_{r-1}^{\prime} p_{r}^{\prime}
\end{array}
\end{aligned}
$$

$$
C=\left(r^{2}-1\right) p_{r-2}^{\prime} p_{r}^{\prime}-r^{2} p_{r-1}^{\prime 2}
$$

Since $a_{1}, a_{2}, \ldots, a_{n-1}$ are not all equal, we have, by the inductive hypothesis,

$$
p_{r-1}^{\prime} p_{r+1}^{\prime}<p_{r}^{\prime 2}, \quad p_{r-2}^{\prime} p_{r}<p_{r-1}^{\prime 2}, \quad p_{r-2}^{\prime} p_{r+1}^{\prime}<p_{r-1}^{\prime} p_{r}^{\prime}
$$

so that

$$
A<-p_{r}^{\prime 2}, \quad B<2 p_{r-1}^{\prime} p_{r}^{\prime}, \quad C<-p_{r-1}^{\prime 2}
$$

and

$$
n^{2}\left(p_{r-1} p_{r+1}-p_{r}^{2}\right)<-\left(p_{r}^{\prime}-a_{n} p_{r-1}^{\prime}\right)^{2} \leqq 0
$$

This proves the theorem. The result is still true when

$$
a_{1}=a_{2}=\ldots=a_{n-1},
$$

because then $a_{n} \neq a_{1}=p_{r}^{\prime} / p_{r-1}^{\prime}$.
It is also possible to prove Theorem 51 by means of identities of the type considered in § 2.21 (1).
54.

$$
p_{r}^{2}-p_{r-1} p_{r+1}=\frac{1}{r(r+1)\binom{n}{r}\binom{n}{r+1}} \sum_{i=0}^{r-1}\binom{2 i}{i} \stackrel{(r, i)}{i+1}
$$

where

$$
(r, i)=\Sigma a_{1}^{2} \ldots a_{r-i-1}^{2} a_{r-i} \ldots a_{r+i-1}\left(a_{r+i}-a_{r+i+1}\right)^{2}
$$

the summation extending over all products formed from the $a$ and of the type shown.
55. $\left\{\frac{n!}{(r-1)!(n-r-1)!}\right\}^{2}\left(p_{r}^{2}-p_{r-1} p_{r+1}\right)=(n-1) \Sigma\left(a_{1}-a_{2}\right)^{2}\left(c_{r-1}^{n-2}\right)^{2}$
$+\frac{2!(n-3)}{(r-1)(n-r-1)} \Sigma\left(a_{1}-a_{2}\right)^{2}\left(a_{3}-a_{4}\right)^{2}\left(c_{r-2}^{n-4}\right)^{2}$
$+\frac{3!(n-5)}{(r-1)(r-2)(n-r-1)(n-r-2)} \Sigma\left(a_{1}-a_{2}\right)^{2}\left(a_{3}-a_{4}\right)^{2}\left(a_{5}-a_{6}\right)^{2}\left(c_{r-3}^{n-6}\right)^{2}$
..+ ,
where $c_{r-1}^{n-2}$ is the sum of the products, $r-1$ at a time, of the $n-2 a$ other than $a_{1}, a_{2}$, and so on.

Theorem 54 is due to Muirhead (1), and Theorem 55 to Jolliffe (1). Theorem 55 gives an 'intuitive' proof of the more general form of Theorem 51 (for real $a$ of either sign) referred to at the beginning of this section.
2.23. A note on definite forms. The identity of Hurwitz and Muirhead proved in § 2.21 (1) shows that, when $a>0$,

$$
a_{1}^{n}+a_{2}^{n}+\ldots+a_{n}^{n}-n a_{1} a_{2} \ldots a_{n}
$$

can be expressed as a sum in which every term is obviously non-negative.

If we write $a_{1}=x_{1}{ }^{2}, a_{2}=x_{2}{ }^{2}, \ldots$, we obtain

$$
\begin{align*}
x_{1}^{2 n} & +\ldots+x_{n}^{2 n}-n x_{1}^{2} x_{2}^{2} \ldots x_{n}^{2}  \tag{2.23.1}\\
& =\frac{1}{2(n-1)!}\left\{\Sigma!\left(x_{1}^{2 n-2}-x_{2}^{2 n-2}\right)\left(x_{1}^{2}-x_{2}^{2}\right)+\ldots\right\}
\end{align*}
$$

Now

$$
\begin{aligned}
\left(x_{1}^{2 n-2}-x_{2}^{2 n-2}\right) & \left(x_{1}^{2}-x_{2}^{2}\right) \\
& =\left(x_{1}^{2}-x_{2}^{2}\right)^{2}\left(x_{1}^{2 n-4}+x_{1}^{2 n-6} x_{2}^{2}+\ldots+x_{2}^{2 n-4}\right)
\end{aligned}
$$

is a sum of squares of polynomials such as $\left(x_{1}{ }^{2}-x_{2}{ }^{2}\right) x_{1}{ }^{n-2}$; and so the right-hand side of (2.23.1) is a sum of squares. Finally, since

$$
\begin{aligned}
x_{1}^{2 n}+ & \ldots+x_{2 n}^{2 n}-2 n x_{1} x_{2} \ldots x_{2 n}=x_{1}^{2 n}+\ldots+x_{n}^{2 n}-n x_{1}^{2} \ldots x_{n}^{2} \\
& +x_{n+1}^{2 n}+\ldots+x_{2 n}^{2 n}-n x_{n+1}^{2} \ldots x_{2 n}^{2}+n\left(x_{1} \ldots x_{n}-x_{n+1} \ldots x_{2 n}\right)^{2}
\end{aligned}
$$

it follows that

$$
\begin{equation*}
F=x_{1}^{2 n}+\ldots+x_{2 n}^{2 n}-2 n x_{1} x_{2} \ldots x_{2 n}=\sum_{i} P_{i}^{2} \tag{2.23.2}
\end{equation*}
$$

where the $P_{i}$ are real polynomials of degree $n$. For example

$$
\begin{aligned}
& x^{6}+y^{6}+z^{6}+u^{6}+v^{6}+w^{6}-6 x y z u v w \\
& \quad=\frac{1}{2}\left(x^{2}+y^{2}+z^{2}\right)\left\{\left(y^{2}-z^{2}\right)^{2}+\left(z^{2}-x^{2}\right)^{2}+\left(x^{2}-y^{2}\right)^{2}\right\} \\
& +\frac{1}{2}\left(u^{2}+v^{2}+w^{2}\right)\left\{\left(v^{2}-w^{2}\right)^{2}+\left(w^{2}-u^{2}\right)^{2}+\left(u^{2}-v^{2}\right)^{2}\right\}+3(x y z-u v w)^{2}
\end{aligned}
$$

is a sum of $9+9+1=19$ squares of real polynomials.
A real form is a homogeneous polynomial $F\left(x_{1}, x_{2}, \ldots, x_{m}\right)$, with real coefficients, in the $m$ real variables $x_{1}, x_{2}, \ldots, x_{m}$. A form $F$ is said to be definite, in a certain region of the variables, if it does not change sign in that region, for example if $F \geqq 0$. We may divide definite forms into positive and negative forms, and it is plainly sufficient to consider positive forms. Thus the form (2.23.2) is
positive in the region of all real values of the variables. It is plain that a form which has this property must be of even degree.

If $F>0$ in a certain region, then $F$ is said to be strictly positive in that region.

The form (2.23.2), and the forms considered in Theorems 7 and 55 (with $x$ for $a$ ), can be expressed as sums of squares of real polynomials, and it is natural to ask whether this is a general property of definite forms. Is it true that, if $F \geqq 0$ for all real $x$, then

$$
F=\sum_{i} P_{i}{ }^{2},
$$

where the $P_{i}$ are real polynomials?
This problem was solved completely by Hilberta. Here we have space only for a few fragmentary remarks. We begin by observing that there are two cases in which the answer is immediate. We denote the degree of $F$ by $2 n$, and the number of variables by $m$.
If $m=2$, so that $F=F(x, y)$ and $n$ is arbitrary, then any real factor $a x+b y$ of $F$ must occur in even multiplicity, and the complex factors must occur in conjugate pairs $a x+b y, \bar{a} x+\bar{b} y$. Hence, grouping the factors appropriately, we obtain

$$
F=p^{2}(q+i r)(q-i r)=(p q)^{2}+(p r)^{2}
$$

where $p, q, r$ are real polynomials.
It is a familiar theorem of algebra ${ }^{b}$ that any definite quadratic form in any number $m$ of variables may be expressed as a sum of at most $m$ squares of real linear forms. Thus the answer is affirmative in the two cases
(1) $m=2, n$ arbitrary,
(2) $m$ arbitrary, $2 n=2$.

Hilbert found a third case
(3) $m=3,2 n=4$,
and proved that any positive biquadratic form in three variables is representable as the sum of three squares of real quadratics. He also proved that in all other cases the answer is negative, there being definite forms of degree $2 n$ in $m$ variables which cannot be represented in the manner proposed.

[^14]Hilbert then suggested the following theorem: any positive $F$ can be expressed as

$$
F=\sum_{i} R_{\imath}{ }^{2},
$$

where $R_{i}$ is a real rational function. An equivalent theorem is: any positive $F$ may be expressed as a quotient of sums of squares of real forms ${ }^{\text {a }}$.

Hilbert ${ }^{\text {b }}$ gave a very difficult proof of these theorems for ternary forms in ( $x, y, z$ ). The general theorems were first proved by Artinc. Artin's proof is very remarkable and comparatively simple, but depends upon the ideas of modern abstract algebra in a manner which makes it impossible for us to reproduce it here.
2.24. A theorem concerning strictly positive forms. The rather fragmentary remarks of $\S 2.23$ form a natural introduction to the simpler problem which we consider here. We are concerned now with forms which are strictly positive in the region of positive $x$. The theorem which we shall prove resembles those of $\S 2.23$ in asserting that a positive form can be represented in a manner which renders its positive character intuitive. It is no longer necessary that the degree of the form should be even.
56.d If the form $\boldsymbol{F}\left(x_{1}, x_{2}, \ldots, x_{m}\right)$ is strictly positive for

$$
x \geqq 0, \quad \Sigma x>0,
$$

then $F$ may be expressed as

$$
F=\frac{G}{H},
$$

where $G$ and $H$ are forms with positive coefficients. In particular we may suppose that
for a suitable $p$.
a It is evident that the first theorem implies the second (with one square only in the denominator). And since

$$
\frac{\Sigma g_{i}{ }^{2}}{\Sigma h_{j}^{2}}=\Sigma \sum_{i, j}\left(\frac{g_{i} h_{j}}{\Sigma h_{j}{ }^{2}}\right)^{2},
$$

the second theorem implies the first.
${ }^{\mathrm{b}}$ Hilbert (2). e Artin (1).
d Pólya (3). The theorem had been proved before (apart from the last clause) by Poincaré (1) when $m=2$ and by Meissner (1) when $m=3$. Meissner's method is applicable in principle in the general case, but does not lead to so simple a result.

For simplicity of writing we suppose $m=3$; no new point of principle arises for general $m$.

The function $F(x, y, z)$ is positive and continuous in the closed region
(2.24.1) $\quad x \geqq 0, \quad y \geqq 0, \quad z \geqq 0, \quad x+y+z=1$,
and has a positive minimum $\mu$ in that region. We write

$$
\begin{equation*}
F(x, y, z)=\Sigma_{n} A_{\alpha \beta \gamma} \frac{x^{\alpha} y^{\beta} z^{\gamma}}{\alpha!\beta!\gamma!}, \tag{2.24.2}
\end{equation*}
$$

the summation being over

$$
\begin{equation*}
\alpha \geqq 0, \quad \beta \geqq 0, \quad \gamma \geqq 0, \quad \alpha+\beta+\gamma=n ; \tag{2.24.3}
\end{equation*}
$$

and
(2.24.4) $\quad \phi(x, y, z ; t)=t^{n} \Sigma_{n} A_{\alpha \beta \gamma}\binom{x t^{-1}}{\alpha}\binom{y t^{-1}}{\beta}\binom{z t^{-1}}{\gamma}$,
where $t>0$ and $\binom{x t^{-1}}{\alpha}, \ldots$ are the usual binomial coefficients, so that

$$
\binom{x t^{-1}}{0}=1
$$

and

$$
t^{\alpha}\binom{x t^{-1}}{\alpha}=\frac{x(x-t)(x-2 t) \ldots\{x-(\alpha-1) t\}}{1.2 .3 \ldots \alpha}
$$

for $\alpha=1,2,3, \ldots$.
It is plain that $\quad \phi(x, y, z ; t) \rightarrow F(x, y, z)$
when $t \rightarrow 0$; and if we write

$$
\phi(x, y, z ; 0)=F(x, y, z),
$$

then $\phi$ is continuous in

$$
x \geqq 0, \quad y \geqq 0, \quad z \geqq 0, \quad x+y+z=1, \quad 0 \leqq t \leqq 1 .
$$

There is therefore an $\epsilon$ such that

$$
\begin{equation*}
\phi(x, y, z ; t)>\phi(x, y, z ; 0)-\frac{1}{2} \mu=F(x, y, z)-\frac{1}{2} \mu \geqq \frac{1}{2} \mu>0 \tag{2.24.5}
\end{equation*}
$$

for $0<t<\epsilon$ and all $x, y, z$ in (2.24.1).
We have also

$$
\begin{equation*}
(x+y+z)^{k-n}=(k-n)!\Sigma_{k-n} \frac{x^{\kappa} y^{\lambda} z^{\mu}}{\kappa!\lambda!\mu!}, \tag{2.24.6}
\end{equation*}
$$

the summation being over

$$
\kappa \geqq 0, \quad \lambda \geqq 0, \quad \mu \geqq 0, \quad \kappa+\lambda+\mu=k-n .
$$

Multiplying (2.24.2) and (2.24.6), we obtain

$$
(x+y+z)^{k-n} F=(k-n)!\Sigma_{n} \Sigma_{k-n} A_{\alpha \beta \gamma} \frac{x^{\alpha+\kappa} y^{\beta+\lambda} z^{\gamma+\mu}}{\alpha!\kappa!\beta!\lambda!\gamma!\mu!} .
$$

Here we write

$$
\alpha+\kappa=a, \quad \beta+\lambda=b, \quad \gamma+\mu=c,
$$

so that $a, b, c$ vary over
(2.24.7) $\quad a \geqq 0, \quad b \geqq 0, \quad c \geqq 0, \quad a+b+c=k$,
and $\alpha, \beta, \gamma$ over
(2.24.8) $\quad 0 \leqq \alpha \leqq a, \quad 0 \leqq \beta \leqq b, \quad 0 \leqq \gamma \leqq c, \quad \alpha+\beta+\gamma=n$.

This gives
(2.24.9)

$$
(x+y+z)^{k-n} F=(k-n)!\Sigma_{k} \frac{x^{a} y^{b} z^{c}}{a!b!c!} \Sigma^{\prime} A_{\alpha \beta \gamma}\binom{a}{\alpha}\binom{b}{\beta}\binom{c}{\gamma} .
$$

In (2.24.9), $\Sigma^{\prime}$ implies summation with respect to $\alpha, \beta, \gamma$ over (2.24.8); but, since $\binom{a}{\alpha}=0,\binom{b}{\beta}=0, \ldots$ if $\alpha>a, \beta>b, \ldots$, we may replace this summation by summation over (2.24.3), i.e. by $\Sigma_{n}$. We thus obtain
(2.24.10)

$$
\begin{aligned}
(x+y+z)^{k-n} F & =(k-n)!\Sigma_{k} \frac{x^{a} y^{b} z^{c}}{a!b!c!} \Sigma_{n} A_{\alpha \beta \gamma}\binom{a}{\alpha}\binom{b}{\beta}\binom{c}{\gamma} \\
& =(k-n)!k^{n} \Sigma_{k} \phi\left(\frac{a}{k}, \frac{b}{k}, \frac{c}{k} ; \frac{1}{k}\right) \frac{x^{a} y^{b} z^{c}}{a!b!c!}
\end{aligned}
$$

The $\phi$ here is positive, by (2.24.5), if $k$ is sufficiently large, and this proves the theorem.
(1) The theorem gives a systematic process for deciding whether a given form $F$ is strictly positive for positive $x$. We multiply repeatedly by $\Sigma x$, and, if the form is positive, we shall sooner or later obtain a form with positive coefficients.

It is instructive to consider the working of the process for

$$
F=x_{1}{ }^{n}+x_{2}{ }^{n}+\ldots+x_{n}{ }^{n}-(n-\epsilon) x_{1} x_{2} \ldots x_{n},
$$

where $\epsilon$ is positive and small. The coefficient of

$$
x_{1}{ }^{i_{1}} x_{2}^{i_{2}} \ldots x_{n}^{i_{n}}
$$

where $i_{1}+i_{2}+\ldots+i_{n}=n(q+1)$, in

$$
\phi=\left(x_{1}+\ldots+x_{n}\right)^{n q} \boldsymbol{F}
$$

is certainly positive if one of the $i$ is 0 . If $i_{1} \geqq 1, \ldots, i_{n} \geqq 1$ it is

$$
\frac{(n q)!}{\left(i_{1}-n\right)!i_{2}!i_{3}!\ldots}+\frac{(n q)!}{i_{1}!\left(i_{2}-n\right)!i_{3}!\ldots}+\ldots-(n-\epsilon) \frac{(n q)!}{\left(i_{1}-1\right)!\left(i_{2}-1\right)!\ldots}
$$

and has the sign of

$$
\begin{aligned}
& \psi\left(i_{1}, i_{2}, \ldots, i_{n}\right)=i_{1}\left(i_{1}-1\right) \ldots\left(i_{1}-n+1\right) \\
& \quad+i_{2}\left(i_{2}-1\right) \ldots\left(i_{2}-n+1\right)+\ldots-(n-\epsilon) i_{1} i_{2} \ldots i_{n} .
\end{aligned}
$$

We require $\psi$ to be positive for all $i_{1}, i_{2}, \ldots, i_{n}$.
If not all of $i_{1}, i_{2}, \ldots$ are equal to $q+1$, there is one, say $i_{1}$, less than $q+1$, and one, say $i_{2}$, greater than $q+1$; and changing $i_{1}, i_{2}$ into $i_{1}+1, i_{2}-1$ changes $\psi$ by

$$
\begin{aligned}
& n\left\{i_{1}\left(i_{1}-1\right) \ldots\left(i_{1}-n+2\right)-\left(i_{2}-1\right)\left(i_{2}-2\right) \ldots\left(i_{2}-n+1\right)\right\} \\
& \quad-(n-\epsilon) i_{3} \ldots i_{n}\left(i_{2}-i_{1}-1\right)<0 .
\end{aligned}
$$

Hence $\psi$ will be positive for all $i$ if it positive when every $i$ is $q+1$. It will be positive in this case if
or

$$
n(q+1) q(q-1) \ldots(q-n+2)>(n-\epsilon)(q+1)^{n}
$$

$$
\left(1-\frac{1}{q+1}\right)\left(1-\frac{2}{q+1}\right) \ldots\left(1-\frac{n-1}{q+1}\right)>1-\frac{\epsilon}{n},
$$

and $a$ fortiori $^{\text {a }}$ if

$$
q+1>\frac{n^{2}(n-1)}{2 \epsilon}
$$

If this condition is satisfied, all the coefficients in $\phi$ are positive.
It follows that $F>0$ for $x>0, \Sigma x>0$. Making $\epsilon \rightarrow 0$, we obtain yet another proof of the theorem of the means, in the form $\Sigma x^{n} \geqq n \Pi x$.
(2) If we write

$$
x_{m}=1-x_{1}-\ldots-x_{m-1}
$$

we obtain a theorem concerning general non-homogeneous polynomials in $m-1$ variables.
57. If a (non-homogeneous) polynomial $f\left(x_{1}, x_{2}, \ldots, x_{m-1}\right)$ is positive in the region $\quad x_{1} \geqq 0, \ldots, x_{m-1} \geqq 0, \quad x_{1}+x_{2}+\ldots+x_{m-1} \leqq 1$, then $f(x)$ can be expressed in the form

$$
f(x)=\Sigma c x_{1}{ }^{a_{1}} \ldots x_{m-1}^{a_{m-1}}\left(1-x_{1}-\ldots-x_{m-1}\right)^{a_{m}},
$$

where the a are non-negative integers and the $c$ are positive.
The theorem is a generalisation of one due to Hausdorff ${ }^{\text {b }}$.

## MISCELLANEOUS THEOREMS AND EXAMPLES ${ }^{\circ}$

58. If $\alpha, \beta, \gamma, \ldots, \lambda$ are greater than -1 , and are all positive or all negative, then $(1+\alpha)(1+\beta) \ldots(1+\lambda)>1+\alpha+\beta+\ldots+\lambda$.
[For the case $\alpha=\beta=\ldots=\lambda$, see James Bernoulli (1, 5, 112).]

[^15]59. If $c>0$ then
$$
|a+b|^{2} \leqq(1+c)|a|^{2}+\left(1+\frac{1}{c}\right)|b|^{2}
$$
for all (real or complex) $a, b$.
[See Bohr (1, 78).]
60. If $\mathfrak{A}_{n}, \mathfrak{G}_{n}$ are the arithmetic and geometric means of $a_{1}, a_{2}, \ldots, a_{n}$, with unit weights, and $\mathfrak{A}_{n+1}, \mathfrak{F}_{n+1}$ those of $a_{1}, a_{2}, \ldots, a_{n}, a_{n+1}$, then
unless
$$
n\left(\mathfrak{A}_{n}-\mathfrak{G}_{n}\right)<(n+1)\left(\mathfrak{A}_{n+1}-\mathfrak{G}_{n+1}\right)
$$
[This theorem, communicated to us by Dr R. Rado, embodies another proof of the theorem of the means. If we write $a_{n+1}=x^{n+1}, \mathfrak{F}_{n}=y^{n+1}$, then the inequality to be proved is
$$
x^{n+1}-(n+1) x y^{n}+n y^{n+1}>0
$$
and this follows from Theorem 41.]
61.
$$
a b \leqq \frac{a^{r}}{r}+\frac{b^{r^{\prime}}}{r^{\prime}} \quad(a>0, b>0, r>1)
$$
with equality when $b=a^{r-1}$.
[Another form of Theorem 37. For this and the next two theorems see Young (1, 5, 6).]
62. $u v \leqq u \frac{u^{p}-1}{p}+\left(\frac{1+p v}{1+p}\right)^{(1+p) / p} \quad(u>0, v>-1 / p, p>0)$.
[Replace $r$ in Theorem 61 by $1+p$ and $a, b$ by $u,(1+p v) /(1+p)$.
63.
$u v \leqq u \log u+e^{v-1} \quad(u>0)$.
[Make $p \rightarrow 0$ in Theorem 62. See also §4.4(5).]
64. If $a>0, a_{1} a_{2} \ldots a_{n}=l^{n}$, then
$$
\left(1+a_{1}\right)\left(1+a_{2}\right) \ldots\left(1+a_{n}\right)>(1+l)^{n}
$$
unless all the $a$ are equal.
[Chrystal (1,51). Example of Theorem 40.]
65. If $a$ and $b$ are positive and $p>1$ or $p<0$, then
$$
\Sigma \frac{a^{p}}{b^{p-1}}>\frac{(\Sigma a)^{p}}{(\Sigma b)^{p-1}}
$$
unless (a) and (b) are proportional. The inequality is reversed if $0<p<1$.
[Radon (1, 1351): transformation of Theorem 13.]
66. If $a>0$ then $\Sigma a \Sigma a^{-1}>n^{2}$, unless all the $a$ are equal.
[From Theorems $\mathbf{7}$ or 9 or 16 or 43.]
67.
$$
\Sigma(a+b) \Sigma \frac{a b}{a+b}<\Sigma a \Sigma b
$$
unless (a) and (b) are proportional.
[Milne (1).]
68. $\Sigma(a+b+c) \Sigma \frac{b c+c a+a b}{a+b+c} \Sigma \frac{a b c}{b c+c a+a b}<\Sigma a \Sigma b \Sigma c$, unless ( $a$ ), (b), (c) are proportional.
69. If $0<r<s$ and
$$
M_{r}(a, b)=\left(\frac{a^{r}+b^{r}}{2}\right)^{1 / r}
$$
then
$$
\Sigma M_{r}(a, b) \Sigma M_{-r}(a, b)<\Sigma M_{s}(a, b) \Sigma M_{-s}(a, b)<\Sigma a \Sigma b,
$$
unless ( $a$ ) and ( $b$ ) are proportional.
70. If $0<k<1$ and
$$
\Sigma a b \geqq A\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}}
$$
for all $b$, then $\Sigma a^{k} \geqq A^{k}$.
[This is an analogue of Theorem 15 for the case $0<k<1$ (when $k^{\prime}<0$ ). If all $a$ are positive, define $b$ by
$$
a b=a^{k}, \quad b=a^{k-1}, \quad b^{k \prime}=a^{k},
$$
when
(i)
$$
\Sigma a^{k} \geqq A\left(\Sigma a^{k}\right)^{1 / k^{\prime}}, \quad \Sigma a^{k} \geqq A^{k} .
$$

If all $a$ vanish, $A$ must be 0 and there is nothing to prove. If some but not all vanish, suppose that $a>0$ in a set $E$ of $\mu$ members and $a=0$ in the complementary set $C E$ of $\nu=n-\mu$ members; and define $b$ as above in $E$, and by $b=G$ in $C E$. Then

$$
\Sigma a^{k}=\Sigma a b \geqq A\left(\sum_{E} a^{k}+\nu G^{k^{\prime}}\right)^{1 / k^{\prime}}=A\left(\Sigma a^{k}+\nu G^{k^{\prime}}\right)^{1 / k^{\prime}}
$$

Making $G \rightarrow \infty$ we obtain (i) again.]
71. If $0<h \leqq a_{\nu} \leqq H, 0<k \leqq b_{\nu} \leqq K$, then

$$
1 \leqq \frac{\Sigma a^{2} \Sigma b^{2}}{(\Sigma a b)^{2}} \leqq \frac{1}{4}\left\{\sqrt{ }\left(\frac{H K}{h k}\right)+\sqrt{ }\left(\frac{h k}{H K}\right)\right\}^{2} .
$$

[See Pólya and Szegö (1, I, 57, 213), where the conditions for equality are given.]
72.

$$
\lim _{r \rightarrow \infty} \frac{\mathfrak{M}_{r+1}^{r+1}(a)}{\mathfrak{M}_{r}^{r}(a)}=\mathfrak{M}_{\infty}(a) .
$$

If all the $a$ are positive, there is a similar theorem for $-\infty$.
73. If

$$
\frac{\mathfrak{A}(a)-\mathfrak{G}(a)}{\mathfrak{A}(a)} \leqq \epsilon<1,
$$

the means being formed with unit weights, then

$$
1+\xi<\frac{a_{\nu}}{\mathfrak{A}}<1+\xi^{\prime},
$$

where $\xi$ is the negative and $\xi^{\prime}$ the positive root of the equation

$$
(1+x) e^{-x}=(1-\epsilon)^{n} .
$$

[See Pólya and Szegö (1, I, 58, 215).]
74. If $\left[\gamma_{1}{ }^{\prime}, \ldots, \gamma_{k}{ }^{\prime}\right] \leqq\left[\gamma_{1}, \ldots, \gamma_{k}\right],\left[\delta_{1}{ }^{\prime}, \ldots, \delta_{l}{ }^{\prime}\right] \leqq\left[\delta_{1}, \ldots, \delta_{l}\right]$
for all values of the variables involved, then

$$
\left[\gamma_{1}^{\prime}, \ldots, \gamma_{k}^{\prime}, \delta_{1}^{\prime}, \ldots, \delta_{l}^{\prime}\right] \leqq\left[\gamma_{1}, \ldots, \gamma_{k}, \delta_{1}, \ldots, \delta_{l}\right]
$$

[It follows from the first hypothesis, and the definition of the sums, that

$$
\left[\gamma_{1}^{\prime}, \ldots, \gamma_{k}^{\prime}, \delta_{1}\right] \leqq\left[\gamma_{1}, \ldots, \gamma_{k}, \delta_{1}\right] ;
$$

and hence, by repetition of the argument, that

$$
\begin{equation*}
\left[\gamma_{1}^{\prime}, \ldots, \gamma_{k}^{\prime}, \delta_{1}, \ldots, \delta_{l}\right] \leqq\left[\gamma_{1}, \ldots, \gamma_{k}, \delta_{1}, \ldots, \delta_{l}\right] . \tag{i}
\end{equation*}
$$

Similarly, using the second hypothesis,

$$
\begin{equation*}
\left[\gamma_{1}^{\prime}, \ldots, \gamma_{k}^{\prime}, \delta_{1}^{\prime}, \ldots, \delta_{l}^{\prime}\right] \leqq\left[\gamma_{1}^{\prime}, \ldots, \gamma_{k}^{\prime}, \delta_{1}, \ldots, \delta_{l}\right] . \tag{ii}
\end{equation*}
$$

The result follows from (i) and (ii).]
75. If $\left(\alpha^{\prime}\right) \prec(\alpha)$ and the $\alpha$ and $\alpha^{\prime}$ are in decreasing order, then there is a greatest non-negative $\delta$ for which

$$
\begin{equation*}
(\beta)=\left(\alpha_{1}^{\prime}+\delta, \alpha_{2}^{\prime}, \ldots, \alpha_{n-1}^{\prime}, \alpha_{n}^{\prime}-\delta\right)<(\alpha) . \tag{i}
\end{equation*}
$$

If $\delta$ has this value then

$$
\begin{equation*}
\left(\alpha_{1}^{\prime}, \alpha_{n}^{\prime}\right) \prec\left(\beta_{1}, \beta_{n}\right) \tag{ii}
\end{equation*}
$$

and
(iii) $\quad\left(\beta_{1}, \ldots, \beta_{k}\right)<\left(\alpha_{1}, \ldots, \alpha_{k}\right),\left(\beta_{k+1}, \ldots, \beta_{n}\right) \prec\left(\alpha_{k+1}, \ldots, \alpha_{n}\right)$
for some $k$ between 1 and $n-1$ inclusive.
[It is plain from the definitions (a) that (i) is true for $\delta=0$, (b) that the set of $\delta$ for which it is true is closed, (c) that, if it is true for a positive $\delta$, then it is true for any smaller positive $\delta$. Hence there is a maximum non-negative $\delta$ for which (i) and (ii) are true.

If $\delta$ has this value then either (a) $\beta_{n}=\alpha_{n}{ }^{\prime}-\delta=0$, or

$$
\text { (b) } \beta_{1}+\ldots+\beta_{k}=\alpha_{1}^{\prime}+\ldots+\alpha_{k}^{\prime}+\delta=\alpha_{1}+\ldots+\alpha_{k}
$$

for a $k<n$; for otherwise we could increase $\delta$ without disturbing (i). In case ( $a$ )

$$
\sum_{1}^{n-1} \alpha_{\nu} \leqq \sum_{1}^{n} \alpha_{\nu}=\sum_{1}^{n} \beta_{\nu}=\sum_{1}^{n-1} \beta_{\nu} \leqq \sum_{1}^{n-1} \alpha_{\nu}
$$

and so $\beta_{1}+\ldots+\beta_{n-1}=\alpha_{1}+\ldots+\alpha_{n-1}$, which is (b) with $k=n-1$. Hence (b) is true in any case; and then (iii) follows from the definitions.

Dr R. Rado, who communicated Theorems 74 and 75 to us, uses them to obtain a new and elegant proof of the sufficiency of Muirhead's criterion (Theorem 45). The result is true for $n=2$ by Lemma 1 of p. 47; let us then suppose that $n>2$, that the conditions (2.18.1), (2.18.2), and (2.18.3) are satisfied, and that the result is true for any number of variables less than $n$. Then, by the inductive hypothesis,
$\left[\alpha_{1}^{\prime}, \alpha_{n}{ }^{\prime}\right] \leqq\left[\beta_{1}, \beta_{n}\right], \quad\left[\beta_{1}, \ldots, \beta_{k}\right] \leqq\left[\alpha_{1}, \ldots, \alpha_{k}\right], \quad\left[\beta_{k+1}, \ldots, \beta_{n}\right] \leqq\left[\alpha_{k+1}, \ldots, \alpha_{n}\right]$.
Hence, using Theorem 74 twice, we obtain

$$
\begin{aligned}
& {\left[\alpha_{1}^{\prime}, \ldots, \alpha_{n}^{\prime}\right]=\left[\alpha_{1}^{\prime}, \beta_{2}, \ldots, \beta_{n-1}, \alpha_{n}{ }^{\prime}\right] \leqq\left[\beta_{1}, \ldots, \beta_{n}\right]} \\
& \left.\quad=\left[\beta_{1}, \ldots, \beta_{k}, \beta_{k+1}, \ldots \beta_{n}\right] \leqq\left[\alpha_{1}, \ldots, \alpha_{k}, \alpha_{k+1}, \ldots \alpha_{n}\right]=\left[\alpha_{1}, \ldots, \alpha_{n}\right] .\right]
\end{aligned}
$$

76. If $a>0$, and $r$ and $s$ are positive integers, then

$$
\Sigma a_{1}^{r} a_{2}^{r} \ldots a_{s}^{r} \Sigma \frac{1}{a_{1}^{s} a_{2}^{s} \ldots a_{r}^{s}}>\binom{n}{r}\binom{n}{s},
$$

unless all the $a$ are equal.
[Generalisation of Theorem 66. By Theorem 45,

$$
[r, r, \ldots(s \text { times }), 0,0, \ldots, 0]>\left[\frac{r s}{n}, \frac{r s}{n}, \ldots, \frac{r s}{n}\right]
$$

Form the corresponding inequality with $r$ and $s$ interchanged and $a$ replaced by $1 / a$, and multiply.]
77. A necessary and sufficient condition that

$$
p_{1}{ }^{a_{1}^{\prime}} p_{2^{a_{2}^{\prime}}} \ldots p_{n^{a_{n}^{\prime}}} \leqq p_{1}^{\alpha_{1}} p_{2}^{\alpha_{2}} \ldots p_{n^{a_{n}}}
$$

where $p_{1}, \ldots$ are defined as in $\S 2.22$, and the $\alpha$ are positive, for all positive $a$, is that

$$
\alpha_{m}^{\prime}+2 \alpha_{m+1}^{\prime}+\ldots+(n-m+1) \alpha_{n}^{\prime} \geqq \alpha_{m}+2 \alpha_{m+1}+\ldots+(n-m+1) \alpha_{n}
$$

for $1 \leqq m \leqq n$, with equality when $m=1$.
[The sufficiency follows from Theorem 51 ; the necessity may be proved on the lines of § 2.19(1). Dougall(1) gives a proof for integral $\alpha$ based on an identity. For certain special cases, such as
see Kritikos (1).]

$$
\begin{gathered}
p_{\mu-\lambda} p_{\mu+\lambda} \leqq p_{\mu-\kappa} p_{\mu+\kappa} \quad(0 \leqq \kappa<\lambda<\mu) \\
p_{\mu_{1}+\mu_{2}+\cdots+\mu_{r}} \leqq p_{\mu_{1}} p_{\mu_{2}} \cdots p_{\mu_{r}}
\end{gathered}
$$

78. The means $\left[\frac{1}{2}, \frac{1}{2}, 0,0, \ldots, 0\right]$ and $\left[\frac{3}{5}, \frac{1}{5}, \frac{1}{5}, 0, \ldots, 0\right]$ are not comparable.
[Example of Theorem 45 and illustration of Theorem 48.]
79. If $a>0$, and $P_{\mu}$ is the arithmetic mean of the $\mu$ th roots of the products of $\mu$ different $a$, then unless all the $a$ are equal.

$$
P_{1}>P_{2}>\ldots>P_{n},
$$

[Smith (1, 440). Example of Theorem 45:

$$
\left.[1,0,0, \ldots, 0]>\left[\frac{1}{2}, \frac{1}{2}, 0, \ldots, 0\right]>\left[\frac{1}{3}, \frac{1}{3}, \frac{1}{3}, \ldots, 0\right]>\ldots\right]
$$

80. If $\mu \geqq 0$, and $x, y, z$ are positive, then

$$
x^{\mu}(x-y)(x-z)+y^{\mu}(y-z)(y-x)+z^{\mu}(z-x)(z-y)>0
$$

unless $x=y=z$.
81. If $v \geqq 0, \delta>0$, and the $a$ are positive and not all equal, then

$$
\left[\nu+2 \delta, 0,0, \alpha_{4}, \ldots\right]-2\left[\nu+\delta, \delta, 0, \alpha_{4}, \ldots\right]+\left[\nu, \delta, \delta, \alpha_{4}, \ldots\right]>0 .
$$

[This result, communicated to us by Prof. I. Schur, is not a consequence of Theorem 45, but follows from Theorem 80, with $\mu=\nu / \delta$.]

## CHAPTER III

## MEAN VALUES WITH AN ARBITRARY FUNCTION AND THE THEORY OF CONVEX FUNCTIONS

3.1. Definitions. The means $\mathfrak{M}_{r}(a)$ and $\mathscr{G}(a)$ are of the form

$$
\begin{equation*}
\mathfrak{M}_{\phi}(a)=\phi^{-1}\{\Sigma q \phi(a)\}, \tag{3.1.1}
\end{equation*}
$$

where $\phi(x)$ is one of the functions

$$
x^{r}, \quad \log x
$$

and $\phi^{-1}(x)$ the inverse function. It is natural to consider more general means of the type (3.1.1), formed with an arbitrary function $\phi$ subject to appropriate conditions. The most obvious conditions to be imposed upon $\phi$ are that it should be continuous and strictly monotonic, in which case it has an inverse $\phi^{-1}$ which satisfies the same conditions.

We require the following preliminary theorem.
82. If (i) $\phi(x)$ is continuous and strictly monotonic in $H \leqq x \leqq K$,
(ii) $H \leqq a_{\nu} \leqq K \quad(\nu=1,2, \ldots, n)$,
(iii) $q_{\nu}>0, \Sigma q_{\nu}=1$,
then (1) there is a unique $\mathfrak{M}$ in $(H, K)$ for which

$$
\begin{equation*}
\phi(\mathfrak{M})=\Sigma q \phi(a), \tag{3.1.2}
\end{equation*}
$$

(2) $\mathfrak{M}$ is greater than some and less than others of the $a$, unless the a are all equal.

Since $\phi(x)$ is continuous and increases or decreases from $\phi(H)$ to $\phi(K)$ when $x$ increases from $H$ to $K$, and $\Sigma q \phi(a)$ lies between these limits, there is just one $\mathfrak{M}$ which satisfies (3.1.2). Also

$$
\Sigma q\{\phi(\mathfrak{M})-\phi(a)\}=0
$$

and some terms must be positive and some negative, unless all are zero. Hence $\mathfrak{M}-a$ is sometimes positive and sometimes negative, unless it is always zero.

We have assumed $\phi(x)$ continuous in the closed interval ( $H, K$ ). The argument is still valid if $\phi(x)$ is continuous and
strictly increasing in $H<x<K$, while $\phi(x) \rightarrow-\infty$ when $x \rightarrow H$, or $\phi(x) \rightarrow+\infty$ when $x \rightarrow K$, provided that we then interpret $\phi(H)$ as meaning $-\infty$, or $\phi(K)$ as meaning $+\infty$, and $\mathfrak{M}$ as being $H$ when $\Sigma q \phi(a)=-\infty$ or $K$ when $\Sigma q \phi(a)=+\infty$. Here $H$ may be $-\infty$ or $K$ may be $+\infty$; a particularly important case is that in which $H=0, K=+\infty$. In the definition which follows, and in all the discussion of the properties of $\mathfrak{M}_{\phi}$ later in this chapter, it is assumed that $\phi$ is strictly monotonic and is either continuous in the closed interval or behaves in the manner just explained.

We write ${ }^{\text {a }}$
(3.1.3) $\mathfrak{M}_{\phi}=\mathfrak{M}_{\phi}(a)=\mathfrak{M} \mathcal{M}_{\phi}(a, q)=\phi^{-1}\{\Sigma q \phi(a)\}=\phi^{-1}\{\mathfrak{U}[\phi(a)]\}$.

The weights $q$ are arbitrary positive numbers whose sum is 1 , and when we compare two means it is to be understood that the weights of the means are the same. For $\phi(x)=x, \log x$ and $x^{r}$, $\mathfrak{M}_{\phi}$ reduces to $\mathfrak{X}$, $\mathfrak{E}$ and $\mathfrak{M}_{r}$ respectively.
3.2. Equivalent means. The mean $\mathbb{M}_{\phi}$ is determined when the function $\phi$ is given. We may ask whether the converse is true: if $\mathfrak{M n}_{\psi}=\mathfrak{M}_{\chi}$ for all $a$ and $q$, is $\psi$ necessarily the same function as $\chi$ ? This question is answered by the theorem which follows.
83. ${ }^{\text {b }}$ In order that

$$
\begin{equation*}
\mathfrak{M}_{\psi}(a)=\mathfrak{M}_{\chi}(a) \tag{3.2.1}
\end{equation*}
$$

for all $a$ and $q$, it is necessary and sufficient that

$$
\begin{equation*}
\chi=\alpha \psi+\beta \tag{3.2.2}
\end{equation*}
$$

where $\alpha$ and $\beta$ are constants and $\alpha \neq 0$.
In the argument which follows we assume $\psi$ and $\chi$ continuous in the closed interval $(H, K)$. It is easy to see that it applies with trivial variations in the exceptional cases mentioned in §3.1. We shall actually prove more than we have stated, viz. that (3.2.2) is a sufficient condition for (3.2.1) to be true for all $a$ and $q$, and that it is a necessary condition for (3.2.1) to be true for all sets of

[^16]two variables and weights. Later (§ 3.7) we shall prove still more, viz. that (3.2.2) is necessary for the truth of (3.2.1) for all sets of two variables and a fixed pair of weights.
(i) If (3.2.2) is satisfied
\[

$$
\begin{aligned}
\chi\left\{\mathfrak{M}_{\chi}(a)\right\} & =\Sigma q \chi(a)=\Sigma q\{\alpha \psi(a)+\beta\} \\
& =\alpha \Sigma q \psi(a)+\beta \\
& =\alpha \psi\left\{\mathfrak{M}_{\psi}(a)\right\}+\beta=\chi\left\{\mathfrak{M}_{\psi}(a)\right\},
\end{aligned}
$$
\]

and so $\mathfrak{M}_{x}=\mathfrak{M}_{\psi}$. Hence the condition is sufficient.
(ii) In proving the condition necessary, we assume only that (3.2.1) is true for all sets of two variables and weights.

In (3.2.1) take

$$
n=2, \quad a_{1}=H, \quad a_{2}=K, \quad q_{1}=\frac{K-t}{K-H}, \quad q_{2}=\frac{t-H}{K-H},
$$

where $H<t<K$. Then
(3.2.3)
$\psi^{-1}\left\{\frac{K-t}{K-H} \psi(H)+\frac{t-H}{K-H} \psi(K)\right\}=\chi^{-1}\left\{\frac{K-t}{K-H} \chi(H)+\frac{t-H}{K-H} \chi(K)\right\}$ for $H<t<K$; and this is also true for $t=H$ and $t=K$. If we denote the common value by $x$ then, as $t$ varies from $H$ to $K$, $x$ assumes all values in $(H, K)$ and

$$
\frac{K-t}{K-H} \psi(H)+\frac{t-H}{K-H} \psi(K)=\psi(x),
$$

$$
\begin{aligned}
\frac{K-t}{K-H} \chi(H)+\frac{t-H}{K-H} \chi(K) & =\frac{\psi(K)-\psi(x)}{\psi(K)-\psi(H)} \chi(H)+\frac{\psi(x)-\psi(H)}{\psi(K)-\psi(H)} \chi(K) \\
& =\alpha \psi(x)+\beta,
\end{aligned}
$$

where $\alpha$ and $\beta$ are independent of $x$. Hence

$$
x=\chi^{-1}\{\alpha \psi(x)+\beta\}
$$

for all $x$ in ( $H, K$ ); and this is (3.2.2). This completes the proof of Theorem 83.

One corollary of Theorem 83 which is sometimes useful is this. Since $-\phi$ is a linear function of $\phi$, and $-\phi$ increases if $\phi$ decreases, we may always suppose, if we please, that the $\phi$ involved in $\mathfrak{M}_{\phi}(x)$ is an increasing function.

Theorem 83 also enables us to elucidate the apparently exceptional status of $\mathfrak{M}_{0}=\mathfrak{G}$ among the means $\mathfrak{M}_{r}$ of Ch. II. Since

$$
\phi_{r}(x)=\int_{1}^{x} t^{r-1} d t=\frac{x^{r}-1}{r}
$$

is a linear function of $x^{r}$, for $r \neq 0$, we have, by Theorem 83,

$$
\mathfrak{M}_{r}(a)=\mathfrak{M}_{\phi_{r}}(a) .
$$

This equation is still valid for $r=0$, since $\phi_{0}(x)=\log x$.
3.3. A characteristic property of the means $\mathfrak{M}_{r}$. It is natural to ask whether there is any simple property of the means of Ch. II which characterises them among the more general means considered here.
84. ${ }^{\text {a }}$ Suppose that $\phi(x)$ is continuous in the open interval $(0, \infty)$, and that

$$
\begin{equation*}
\mathfrak{M}_{\phi}(k a)=k \mathfrak{M}_{\phi}(a) \tag{3.3.1}
\end{equation*}
$$

for all positive $a, q$, and $k$. Then $\mathfrak{M}_{\phi}(a)$ is $\mathfrak{M}_{r}(a)$. In other words, the means $\mathfrak{M r}_{r}$ are the only homogeneous means $\mathfrak{M}_{\phi}$.

Naturally (3.3.1) does not imply $\phi=x^{r}$ (or $\log x$ ); for, by Theorem 83, we can replace $\phi$ by $\alpha \phi+\beta$ without changing $\mathfrak{M}_{\phi}$.

That (3.3.1) is true when $\phi=x^{r}$ or $\phi=\log x$ is obvious. We now assume (3.3.1) and deduce the form of $\phi$. After Theorem 83, we may suppose that

$$
\begin{equation*}
\phi(1)=0 ; \tag{3.3.2}
\end{equation*}
$$

for we may replace $\phi(x)$ by $\phi(x)-\phi(1)$.
We write (3.3.1) in the form

$$
\mathfrak{M}_{\phi}(a)=k^{-1} \mathfrak{M}_{\phi}(k a)=k^{-1} \phi^{-1}\{\Sigma q \phi(k a)\}=\mathfrak{M}_{\psi}(a),
$$

where

$$
\psi(x)=\phi(k x) .
$$

It follows from Theorem $83{ }^{b}$ that

$$
\begin{equation*}
\phi(k x)=\alpha(k) \phi(x)+\beta(k), \tag{3.3.3}
\end{equation*}
$$

where $\alpha(k)$ and $\beta(k)$ are functions of $k$, and $\alpha(k) \neq 0$; and from (3.3.2) and (3.3.3) that

$$
\begin{equation*}
\phi(k)=\beta(k) . \tag{3.3.4}
\end{equation*}
$$

a Nagumo (1), de Finetti (1), Jessen (4). The following simple version of de Finetti's proof was communicated to us by Dr Jessen.
b If we used one of the more precise forms of Theorem 83, referred to after its enunciation in §3.2, we should obtain a more precise form of Theorem 84, in which homogeneity was only assumed for restricted classes of variables or weights.

If we substitute from (3.3.4) into (3.3.3), and write $y$ for $k$, we find that

$$
\begin{equation*}
\phi(x y)=\alpha(y) \phi(x)+\phi(y) \tag{3.3.5}
\end{equation*}
$$

for all positive $x$ and $y$.
Similarly

$$
\begin{equation*}
\phi(x y)=\alpha(x) \phi(y)+\phi(x) ; \tag{3.3.6}
\end{equation*}
$$

and (3.3.5) and (3.3.6) give

$$
\frac{\alpha(x)-1}{\phi(x)}=\frac{\alpha(y)-1}{\phi(y)} .
$$

Each of these functions must reduce to a constant $c$, so that $\alpha(y)=1+c \phi(y)$. It then follows from (3.3.5) that
(3.3.7) $\quad \phi(x y)=c \phi(x) \phi(y)+\phi(x)+\phi(y)$.

In discussing this functional equation, we must distinguish two cases.
(1) If $c=0,(3.3 .7)$ reduces to the classical equation

$$
\phi(x y)=\phi(x)+\phi(y) .
$$

The most general solution, continuous for $x>0$, is ${ }^{\mathrm{b}} \phi=C \log x$.
(2) If $c \neq 0$, we put $\quad c \phi(x)+1=f(x)$,
and the equation reduces to

$$
f(x y)=f(x) f(y),
$$

whose general solution is $f=x^{r}$. Hence

$$
\phi(x)=\frac{x^{r}-1}{c} .
$$

3.4. Comparability. Our general remarks on the 'comparability' of functions of the $a$ (§1.6) suggest the following problem: given two functions $\psi$ and $\chi$, each continuous and strictly monotonic in ( $H, K$ ), are $\mathfrak{M}_{\psi}$ and $\mathfrak{M}_{\chi}$ comparable; i.e. is there an inequality

$$
\begin{equation*}
\mathfrak{M}_{\psi} \leqq \mathfrak{M}_{\chi} \tag{3.4.1}
\end{equation*}
$$

[^17](or one in the opposite sense) valid for all $a$ and $q$ ? Theorem 16 tells us that the answer is affirmative when $\psi$ and $\chi$ are powers.

We write $\quad \chi\left\{\psi^{-1}(x)\right\}=\phi(x)$.
Then $\phi$ is continuous and strictly monotonic, and has the inverse $\phi^{-1}=\psi \chi^{-1}$. We also write

$$
x=\psi(a), \quad a=\psi^{-1}(x)
$$

Then the $x$ are arbitrary numbers between $\psi(H)$ and $\psi(K)$; and (3.4.1) takes the form

$$
\begin{equation*}
\phi(\Sigma q x) \leqq \Sigma q \phi(x) \tag{3.4.2}
\end{equation*}
$$

(for all $q$ ) if $\chi$ is increasing, the reversed form if $\chi$ is decreasing.
We thus obtain
85. If $\psi$ and $\chi$ are continuous and strictly monotonic, then a necessary and sufficient condition that $\mathfrak{M}_{\psi}$ and $\mathfrak{M}_{\chi}$ should be comparable is that $\phi=\chi \psi^{-1}$ should satisfy (3.4.2), or the reversed inequality.

In what follows, we examine this class of functions $\phi$ in detail. For arbitrary weights $p$, (3.4.2) becomes

$$
\begin{equation*}
\phi\left(\frac{\Sigma p x}{\Sigma p}\right) \leqq \frac{\Sigma p \phi(x)}{\Sigma p} \tag{3.4.3}
\end{equation*}
$$

3.5. Convex functions. The function $\phi$ of $\S 3.4$ was the resultant of two monotonic functions, and therefore itself monotonic; but now we consider a $\phi$ subject to (3.4.2) only.

The simplest case of (3.4.2) is

$$
\begin{equation*}
\phi\left(\frac{x+y}{2}\right) \leqq \frac{\phi(x)+\phi(y)}{2} \tag{3.5.1}
\end{equation*}
$$

A function which satisfies (3.4.2) satisfies (3.5.1), but the class of functions satisfying (3.5.1) is more general. We shall however show that the two inequalities are equivalent for functions subject to certain not very restrictive conditions.

A function which satisfies (3.5.1) in a certain interval is called convex in that interval. If $-\phi$ is convex, $\phi$ is concave. We may also define convexity or concavity in an open interval. It is often convenient to admit infinite values at the ends of the in-
terval; it is obvious that such values must be positive for convex and negative for concave functions.

The foundations of the theory of convex functions are due to Jensen (2). ${ }^{\text {a }}$ Geometrically, (3.5.1) means that the middle point of any chord of the curve $y=\phi(x)$ lies above or on the curve; here a curve means any, not necessarily continuous, graph. The inequality
(3.5.2) $\quad \phi\left(q_{1} x_{1}+q_{2} x_{2}\right) \leqq q_{1} \phi\left(x_{1}\right)+q_{2} \phi\left(x_{2}\right)$
(for all $q$ ) asserts that the whole chord lies above or on the curve, and the general inequality (3.4.2) asserts that the centre of gravity of any number of arbitrarily weighted points of the curve lies above or on the curve. It is geometrically intuitive that, when the curve is continuous, the weakest condition implies the stronger, and we shall find that much more than this is confirmed by our analysis. We might have taken (3.4.2) or (3.5.2) as our definition of convexity, but we have followed Jensen in starting from the weakest definition. The most natural definitions are perhaps (3.5.2) and another which we discuss in § 3.19. There is some logical interest in assuming as little as possible.

It is sometimes useful to have a definition of the convexity or concavity of a finite or enumerably infinite set of numbers. We shall say that the set $a_{1}, \ldots, a_{n}$ is convex if

$$
2 a_{\nu} \leqq a_{\nu-1}+a_{\nu+1} \quad(\nu=2,3, \ldots, n-1),
$$

i.e. if the second differences of the set are non-negative.

Thus we can state Theorem 51, in the less exact form with ' $\leqq$ ', by saying that the set $\log p$ is concave; the full theorem is that $\log p$ is strictly concave (see § 3.8) unless the $a$ are equal. When two products of powers of the $p$ are comparable, the inequality which holds between them may be deduced (substantially as Theorem 52 was deduced from Theorem 51) from the concavity of $\log p$. This is the kernel of Theorem 77.
3.6. Continuous convex functions. We now proceed to investigate the simplest case in which (3.4.2) and (3.5.1) are equivalent.

If $\phi(x)$ satisfies (3.5.1), we have

$$
\begin{gathered}
4 \phi\left(\frac{x_{1}+x_{2}+x_{3}+x_{4}}{4}\right) \leqq 2 \phi\left(\frac{x_{1}+x_{2}}{2}\right)+2 \phi\left(\frac{x_{3}+x_{4}}{2}\right) \\
\leqq \phi\left(x_{1}\right)+\phi\left(x_{2}\right)+\phi\left(x_{3}\right)+\phi\left(x_{4}\right),
\end{gathered}
$$

a Though Hölder (1) had considered the inequality (3.4.2) before Jensen.
and so on. We thus prove

$$
\begin{equation*}
\phi\left(\frac{x_{1}+x_{2}+\ldots+x_{n}}{n}\right) \leqq \frac{\phi\left(x_{1}\right)+\phi\left(x_{2}\right)+\ldots+\phi\left(x_{n}\right)}{n} \tag{3.6.1}
\end{equation*}
$$

for a particular sequence of values of $n$, viz. $n=2^{m}$.
To prove (3.6.1) true generally, it is enough to prove also that, if it is true for $n$, it is true for $n-1 .{ }^{\text {a }}$ Suppose then that (3.6.1) has been proved for $n$ numbers, and that $x_{1}, x_{2}, \ldots, x_{n-1}$ are given. Taking $x_{n}$ to be the arithmetic mean $\mathfrak{H}$ (with equal weights) of the $n-1$ numbers, and applying (3.6.1), we obtain

$$
\begin{aligned}
\phi(\mathfrak{A})= & \phi\left\{\frac{(n-1) \mathfrak{A}+\mathfrak{A}}{n}\right\}=\phi\left(\frac{x_{1}+x_{2}+\ldots+x_{n-1}+\mathfrak{H}}{n}\right) \\
& \leqq \frac{\phi\left(x_{1}\right)+\phi\left(x_{2}\right)+\ldots+\phi\left(x_{n-1}\right)+\phi(\mathfrak{H})}{n},
\end{aligned}
$$

and so

$$
\phi(\mathfrak{U}) \leqq \frac{\phi\left(x_{1}\right)+\phi\left(x_{2}\right)+\ldots+\phi\left(x_{n-1}\right)}{n-1}
$$

which is (3.6.1) with $n-1$ for $n$. Hence (3.6.1) is true generally.
Next, supposing that, in (3.6.1), the $x$ form appropriate groups of equal numbers, we obtain (3.4.2) for any commensurable $q$.

Finally, if $\phi(x)$ is continuous, we can prove (3.4.2) without restriction on the $q$; for we may replace the $q$ by commensurable approximations and proceed to the limit. We thus obtain
86. Any continuous convex function satisfies (3.4.2).

As an application, we may consider Theorem 17. If $s=\frac{1}{2}(r+t)$ we have, by Theorem 7,

$$
\left(\Sigma p a^{8}\right)^{2} \leqq \Sigma p a^{r} \Sigma p a^{t}
$$

or

$$
\left\{\mathfrak{M}_{s}^{s}(a)\right\}^{2} \leqq \mathfrak{M}_{r}^{r}(a) \mathfrak{M}_{t}^{t}(a),
$$

or

$$
\log \mathfrak{M}_{s}^{s}(a) \leqq \frac{1}{2}\left\{\log \mathfrak{M}_{r}{ }^{r}(a)+\log \mathfrak{M}_{t}{ }^{t}(a)\right\} .
$$

In other words
87. $\log \mathfrak{M}_{r}{ }^{r}(a)=r \log \mathfrak{M}_{r}(a)$ is a convex function of $r$.

From this, by appealing to Theorem 86 (or repeating the argument by which this theorem was proved), we deduce Theorem 17 (apart from the specification of the cases of equality).

[^18]3.7. An alternative definition ${ }^{\mathrm{a}}$. We characterised a convex function $\phi(x)$ in $\S 3.5$ by the fact that the middle point of a chord of the curve $y=\phi(x)$ lies above or on the curve. Riesz and Jessen have made an observation which is interesting and sometimes important in applications ${ }^{\text {b }}$, viz. that, when $\phi(x)$ is continuous, it is sufficient to require that some point of the chord should lie above or on the curve.
88. If $\phi(x)$ is continuous, and there is at least one point of every chord of the curve $y=\phi(x)$, besides the end points of the chord, which lies above or on the curve, then every point of every chord lies above or on the curve, so that $\phi(x)$ is convex.

Suppose that $P Q$ is a chord, and $R$ a point on the chord below the curve. Then there is a last point $S$ on $P R$ and a first point $T$ on $R Q$ in which the curve meets the chord: $S$ may be $P$ and $T$ may be $Q$. The chord $S T$ lies entirely below the curve, contradicting the hypothesis.

This remark gives us an alternative proof of Theorem 86. If $\phi(x)$ is convex, the middle point of any chord lies above or on the curve. Hence, as we have proved, every point of the chord lies above or on the curve. That is to say

$$
\phi\left(q_{1} x_{1}+q_{2} x_{2}\right) \leqq q_{1} \phi\left(x_{1}\right)+q_{2} \phi\left(x_{2}\right)
$$

if $q_{1}>0, q_{2}>0, q_{1}+q_{2}=1$, but $q_{1}$ and $q_{2}$ are otherwise arbitrary. We may then proceed by induction. If $q_{1}+q_{2}+q_{3}=1$, then

$$
\begin{align*}
& \phi\left(q_{1} x_{1}+q_{2} x_{2}+q_{3} x_{3}\right)=\phi\left\{q_{1} x_{1}+\left(q_{2}+q_{3}\right) \frac{q_{2} x_{2}+q_{3} x_{3}}{q_{2}+q_{3}}\right\}  \tag{3.7.1}\\
& \quad \leqq q_{1} \phi\left(x_{1}\right)+\left(q_{2}+q_{3}\right) \phi\left(\frac{q_{2} x_{2}+q_{3} x_{3}}{q_{2}+q_{3}}\right) \\
& \quad \leqq q_{1} \phi\left(x_{1}\right)+\left(q_{2}+q_{3}\right) \frac{q_{2} \phi\left(x_{2}\right)+q_{3} \phi\left(x_{3}\right)}{q_{2}+q_{3}} \\
& \quad=q_{1} \phi\left(x_{1}\right)+q_{2} \phi\left(x_{2}\right)+q_{3} \phi\left(x_{3}\right)
\end{align*}
$$

and so generally.
A corollary of Theorem 88 is
89. If $\phi(x)$ is continuous, and every chord of $y=\phi(x)$ meets the curve in a point distinct from its end-points, then $\phi(x)$ is linear.

[^19]By Theorem 88, every point of every chord lies on or above the curve. But Theorem 88 remains true if 'above' is replaced by 'below' in hypothesis and conclusion. Hence every chord of the curve coincides with the curve.

From Theorem 89 we can deduce the refinement on Theorem 83 referred to in §3.2. Suppose that

$$
\psi^{-1}\left\{q_{1} \psi\left(a_{1}\right)+q_{2} \psi\left(a_{2}\right)\right\}=\chi^{-1}\left\{q_{1} \chi\left(a_{1}\right)+q_{2} \chi\left(a_{2}\right)\right\}
$$

for fixed $q_{1}, q_{2}$ and arbitrary $a$. Writing $\chi \psi^{-1}=\phi, \psi(a)=x$, $a=\psi^{-1}(x)$, we obtain

$$
\phi\left(q_{1} x_{1}+q_{2} x_{2}\right)=q_{1} \phi\left(x_{1}\right)+q_{2} \phi\left(x_{2}\right),
$$

so that one point at any rate of every chord of $y=\phi(x)$ lies on the curve. It follows from Theorem 89 that $\phi$ is linear.
3.8. Equality in the fundamental inequalities. We now suppose $\phi(x)$ continuous and convex, and consider when equality can occur in (3.5.1), (3.5.2), or (3.4.2).

Suppose that $x_{1}<x_{3}<x_{2}$, that $x_{3}=q_{1} x_{1}+q_{2} x_{2}$, and that $P_{1}$, $P_{2}, \ldots$ are the points on the curve $y=\phi(x)$ corresponding to $x_{1}, x_{2}, \ldots$ If $\phi(x)$ is not linear in $\left(x_{1}, x_{2}\right)$, there is an $x_{4}$ in $\left(x_{1}, x_{2}\right)$ such that $P_{4}$ lies below the line $P_{1} P_{2}$. Suppose for example that $x_{4}$ lies in $\left(x_{1}, x_{3}\right)$. Then $x_{3}$ lies in $\left(x_{4}, x_{2}\right)$, and $P_{3}$ lies on or below $P_{4} P_{2}$, and therefore below $P_{1} P_{2}$. Hence (3.5.2) holds with inequality. It follows that equality can occur in (3.5.2) only when $\phi(x)$ is linear in $\left(x_{1}, x_{2}\right)$.

This conclusion is easily extended to the general inequality (3.4.2). Suppose, for example, that there is equality when $n=3$, and that $x_{1}<x_{2}<x_{3}$. Then all the signs of inequality in (3.7.1) must reduce to equality, and $\phi(x)$ must be linear in each of the intervals

$$
\left(x_{1}, \frac{q_{2} x_{2}+q_{3} x_{3}}{q_{2}+q_{3}}\right), \quad\left(x_{2}, x_{3}\right)
$$

and therefore over ( $x_{1}, x_{3}$ ).
We have thus proved
90. If $\phi(x)$ is continuous and convex, then

$$
\begin{align*}
& \phi(\Sigma q x)<\Sigma q \phi(x)  \tag{3.8.1}\\
& \phi\left(\frac{\Sigma p x}{\Sigma p}\right)<\frac{\Sigma p \phi(x)}{\Sigma p} \tag{3.8.2}
\end{align*}
$$

unless either (i) all the $x$ are equal or (ii) $\phi(x)$ is linear in an interval including all the $x$.
91. Any chord of a continuous convex curve lies entirely above the curve, except for its end-points, or coincides with it.

We may say that $\phi(x)$ is strictly convex if

$$
\begin{equation*}
\phi\left(\frac{x+y}{2}\right)<\frac{1}{2}\{\phi(x)+\phi(y)\} \tag{3.8.3}
\end{equation*}
$$

for every unequal pair $x, y$. Since a strictly convex function cannot be linear in any interval, any such function, if continuous, satisfies (3.8.1) and (3.8.2), unless all the $x$ are equal.
3.9. Restatements and extensions of Theorem 85.a We may restate Theorem 85 in the form
92. If $\psi$ and $\chi$ are continuous and strictly monotonic, and $\chi$ is increasing, then a necessary and sufficient condition that $\mathfrak{M}_{\psi} \leqq \mathbb{M}_{\chi}$ for all $a$ and $q$ is that $\phi=\chi \psi^{-1}$ should be convex ${ }^{b}$.

We shall say in these circumstances that $\chi$ is convex with respect to $\psi$. Thus $t^{s}$ is convex with respect to $t^{r}$ when $s \geqq r>0$.

The curve $y=\phi(x)$ has the parametric representation

$$
x=\psi(t), \quad y=\chi(t)
$$

The chord through the points on the curve corresponding to $t=t_{1}$ and $t=t_{2}$ is

$$
x=\psi(t), \quad y=\psi^{*}(t)
$$

where $\quad \psi^{*}(t)=\frac{\psi\left(t_{2}\right)-\psi(t)}{\psi\left(t_{2}\right)-\psi\left(t_{1}\right)} \chi\left(t_{1}\right)+\frac{\psi(t)-\psi\left(t_{1}\right)}{\psi\left(t_{2}\right)-\psi\left(t_{1}\right)} \chi\left(t_{2}\right)$
is the function

$$
\alpha \psi(t)+\beta
$$

which assumes the values $\chi\left(t_{1}\right)$ and $\chi\left(t_{2}\right)$ for $t=t_{1}$ and $t=t_{2}$. We may call $y=\psi^{*}(x)$ the $\psi$-chord of $y=\chi(x)$. In order that $\chi$ should be convex with respect to $\psi$, it is necessary and sufficient that $\chi \leqq \psi^{*}$, i.e. that every point of any $\psi$-chord of $\chi$ should lie on or above the curve.

[^20]Theorem 92 may be generalised as follows. Suppose that

$$
a=a_{\nu_{1}, \nu_{2}}, \ldots, \nu_{m}
$$

is a function of $m$ variables $\nu_{1}, \nu_{2}, \ldots, \nu_{m}$, and that

$$
\mathfrak{M}_{\psi_{m}}^{\nu_{m}} \ldots \mathfrak{m}_{\psi_{2}}^{\nu_{2}} \mathfrak{M}_{\psi_{1}}^{\nu_{1}}(a)
$$

is the result of taking means with respect to $\nu_{1}, \nu_{2}, \ldots \nu_{m}$ in succession.
93. Suppose that $\psi_{\mu}$ and $\chi_{\mu}$ are continuous and strictly monotonic, and that $\chi_{\mu}$ is increasing. Then, in order that

$$
\mathfrak{M}_{\psi_{m}}^{\nu_{m}} \ldots \mathfrak{M}_{\psi_{1}}^{\nu_{1}}(a) \leqq \mathfrak{M}_{\chi_{m}}^{\nu_{m}} \ldots \mathfrak{M}_{\chi_{1}}^{\nu_{1}}(a)
$$

for all $a$ and $q$, it is necessary and sufficient that every $\chi_{\mu}$ should be convex with respect to the corresponding $\psi_{\mu}$.

It is understood, of course, that the weights involved in the operations $\mathfrak{M}_{\psi_{\mu}}$ and $\mathfrak{M}_{\chi_{\mu}}$ are the same, though they will generally vary with $\mu$. That the conditions are sufficient follows at once from Theorem 92. To see that they are necessary, we have only to suppose $a$ to be a function of a single $\nu_{\mu}$.
3.10. Twice differentiable convex functions. We postpone to §3.18 any further discussion of the general properties of convex functions, and consider now a particularly important sub-class of such functions, viz. those which possess a second differential coefficient.
94. Suppose that $\phi(x)$ possesses a second derivative $\phi^{\prime \prime}(x)$ in the open interval ( $H, K$ ). Then a necessary and sufficient condition that $\phi(x)$ should be convex in the interval is

$$
\begin{equation*}
\phi^{\prime \prime}(x) \geqq 0 .^{\mathrm{a}} \tag{3.10.1}
\end{equation*}
$$

(i) The condition is necessary. Replacing $\frac{1}{2}(x+y)$ and $\frac{1}{2}(x-y)$ by $t$ and $h$ in (3.5.1), and supposing that $x>y$, so that $h>0$, we obtain

[^21]\[

$$
\begin{equation*}
\phi(t+h)+\phi(t-h)-2 \phi(t) \geqq 0 \tag{3.10.2}
\end{equation*}
$$

\]

for all $t, h$ such that the arguments lie in the interval.
Suppose now that $\phi^{\prime \prime}(t)<0$. Then there are positive numbers $\delta$ and $h$ such that

$$
\phi^{\prime}(t+u)-\phi^{\prime}(t-u)<-\delta u
$$

for $0<u \leqq h$. Integrating this inequality from $u=0$ to $u=h$, we obtain

$$
\phi(t+h)+\phi(t-h)-2 \phi(t)<-\frac{1}{2} \delta h^{2},
$$

in contradiction with (3.10.2).
(ii) The condition is sufficient. We prove that $\phi$ satisfies (3.4.2). In fact, if $X=\Sigma q x$, we have

$$
\phi\left(x_{\nu}\right)=\phi(X)+\left(x_{\nu}-X\right) \phi^{\prime}(X)+\frac{1}{2}\left(x_{\nu}-X\right)^{2} \phi^{\prime \prime}\left(\xi_{\nu}\right)
$$

for some $\xi_{\nu}$ between $X$ and $x_{\nu}$, and so

$$
\Sigma q \phi(x) \geqq \phi(X)=\phi(\Sigma q x)
$$

If $\phi^{\prime \prime}(x)>0$, there can be equality only if every $x$ is equal to $X$. We have therefore proved
95. ${ }^{\text {a }}$ If $\phi^{\prime \prime}(x)>0$, then $\phi(x)$ is strictly convex and satisfies (3.8.1) and (3.8.2), unless all the $x$ are equal.
3.11. Applications of the properties of twice differentiable convex functions. The following theorem, which follows from Theorems 95 and $85,{ }^{\text {b }}$ will be found particularly useful in applications.
96. If $\psi$ and $\chi$ are monotonic, $\chi$ increasing, $\phi=\chi \psi^{-1}$, and $\phi^{\prime \prime}>0$, then $\mathfrak{M}_{\psi}<\mathfrak{M}_{\chi}$, unless all the a are equal.

Examples. (1) If $\psi=\log x, \chi=x$, then $\phi=\chi \psi^{-1}=e^{x}$. Theorem 96 reduces to Theorem 9.
(2) If $\psi=x^{r}, \chi=x^{s}$, where $0<r<s$, then $\phi=x^{s / r}, \phi^{\prime \prime}>0$. Theorem 96 gives Theorem 16 (for positive indices). The other cases of Theorem 16 may be derived similarly.
(3) Suppose that $\phi=x^{k}$, where $k$ is not 0 or 1 . Then $\phi$ is convex in $(0, \infty)$ if $k<0$ or $k>1$, concave if $0<k<1$. Supposing $k>1$, and
a Hölder (1).
b More strictly, from Theorem 95 and the proof of Theorem 85. By Theorem 95, (3.4.2) is true with inequality, and so $\mathfrak{M}_{\psi}<\mathfrak{M}_{\chi}$, unless the $a$ are equal. applying Theorem 95, we find

$$
\Sigma q x<\left(\Sigma q x^{k}\right)^{1 / k}
$$

or

$$
\Sigma p x<\left(\Sigma p x^{k}\right)^{1 / k}(\Sigma p)^{1 / k^{\prime}},
$$

unless all the $x$ are equal. If we write $p x=a b, p x^{k}=a^{k}$, we obtain Theorem 13, for $k>1$. The other cases follow similarly.
(4) Suppose that $\phi=\log \left(1+e^{x}\right)$, so that

$$
\dot{\phi}^{\prime \prime}(x)=\frac{e^{x}}{\left(1+e^{x}\right)^{2}}>0
$$

and that the abscissae and weights in (3.8.1) are $\log \left(a_{2} / a_{1}\right)$, $\log \left(b_{2} / b_{1}\right), \ldots$, and $\alpha, \beta, \ldots$. We obtain

$$
a_{1}^{\alpha} b_{1}^{\beta} \ldots l_{1}^{\lambda}+a_{2}^{\alpha} b_{2}^{\beta} \ldots l_{2}^{\lambda}<\left(a_{1}+a_{2}\right)^{\alpha}\left(b_{1}+b_{2}\right)^{\beta} \ldots\left(l_{1}+l_{2}\right)^{\lambda},
$$

unless $a_{2} / a_{1}=b_{2} / b_{1}=\ldots$ (Theorem 40: H for any number of sets of two numbers).
(5) Suppose that $\phi=\left(1+x^{r}\right)^{1 / r}$, where $r$ is not 0 or 1 , and that the abscissae and weights in (3.8.2) are $a_{2} / a_{1}, b_{2} / b_{1}, \ldots$ and $a_{1}, b_{1}, \ldots$. In this case $\phi$ is convex if $r>1$ and concave if $r<1$. We find, for example, that

$$
\begin{aligned}
\left\{\left(a_{1}+b_{1}+\ldots+l_{1}\right)^{r}\right. & \left.+\left(a_{2}+b_{2}+\ldots+l_{2}\right)^{r}\right\}^{1 / r} \\
& <\left(a_{1}^{r}+a_{2}^{r}\right)^{)^{\prime / r}}+\left(b_{1}^{r}+b_{2}^{r}\right)^{1 / r}+\ldots+\left(l_{1}^{r}+l_{2}^{r}\right)^{1 / r}
\end{aligned}
$$

if $r>1$ and $a_{2} / a_{1}, b_{2} / b_{1}, \ldots$ are not all equal (M for any number of sets of two numbers). It will be remembered that both H and M can be extended to sets of more numbers inductively.
(6) 97. If $a>0, p>0$, then

$$
\exp \left(\frac{\Sigma p \log a}{\Sigma p}\right)<\frac{\Sigma p a}{\Sigma p}<\exp \left(\frac{\Sigma p a \log a}{\Sigma p a}\right)
$$

unless all the a are equal.
We write this with $p$ instead of $q$ for the sake of symmetry. The first inequality is (3.8.2) reversed, with $\phi(x)=\log x$, a concave function. It is equivalent to G (Theorem 9). The second inequality is (3.8.2), with $\phi(x)=x \log x$, a convex function.
3.12. Convex functions of several variables. Suppose that $D$ is a convex domain in the plane of $(x, y)$, that is to say, a domain
which contains the whole of the segment of a straight line which connects any two of its points ${ }^{\text {a }}$. A function $\Phi(x, y)$ is said to be convex in $D$ if it is defined everywhere in $D$ and

$$
\begin{equation*}
\Phi\left(\frac{x_{1}+x_{2}}{2}, \frac{y_{1}+y_{2}}{2}\right) \leqq \frac{1}{2}\left\{\Phi\left(x_{1}, y_{1}\right)+\Phi\left(x_{2}, y_{2}\right)\right\} \tag{3.12.1}
\end{equation*}
$$

for all $\left(x_{1}, y_{1}\right)$ and ( $x_{2}, y_{2}$ ) of $D .{ }^{\text {b }}$ The definition asserts more than convexity in $x$ and $y$ separately; thus $x y$ is a convex function of $x$ for every $y$, and a convex function of $y$ for every $x$, but it is not a convex function of $x$ and $y$.

It is often convenient to use an alternative form of the definition just given. Suppose that $x, y, u, v$ are given, and consider the values of $t$ (if any) for which ( $x+u t, y+v t$ ) belongs to $D$. Since $D$ is convex, these values form an interval (which may be nul). Then we say that $\Phi(x, y)$ is convex in $D$ if

$$
\begin{equation*}
\chi(t)=\Phi(x+u t, y+v t) \tag{3.12.2}
\end{equation*}
$$

is, for every $x, y, u, v$, a convex function of $t$ throughout the interval of $t$ in question. The definition is equivalent to that which we gave before, since, if

$$
x+u t_{1}=x_{1}, y+v t_{1}=y_{1}, \quad x+u t_{2}=x_{2}, y+v t_{2}=y_{2}
$$

(3.12.1) becomes

$$
\chi\left(\frac{t_{1}+t_{2}}{2}\right) \leqq \frac{1}{2}\left\{\chi\left(t_{1}\right)+\chi\left(t_{2}\right)\right\} .
$$

$\Phi$ is said to be concave if $-\Phi$ is convex.
If $z=\Phi(x, y)$ is the equation of a surface in rectangular Cartesian coordinates, (3.12.1) asserts that the middle point of any chord of the surfacelies above or on the corresponding point of the

[^22]surface. If the surface is continuous we can deduce that the whole chord lies above or on the surface, and that the centre of gravity of any number of arbitrarily weighted points of the surface lies above or on the surface. This is what is asserted in the following theorem.
98. If $\Phi(x, y)$ is convex and continuous, then
\[

$$
\begin{equation*}
\Phi(\Sigma q x, \Sigma q y) \leqq \Sigma q \Phi(x, y) \tag{3.12.3}
\end{equation*}
$$

\]

The proof is the same as that of Theorem 86, except for the obvious changes of notation.

There is also a theorem corresponding to Theorem 88; it is sufficient to assert that no chord of the surface lies (except for its end-points) entirely below the surface. All the other remarks of § 3.7 remain true with the obvious changes.

A theorem corresponding to Theorems 94 and 95 is
99. If $\Phi(x, y)$ is twice differentiable in an open $D$, then a necessary and sufficient condition that it should be convex in $D$ is that the quadratic form

$$
Q=\Phi_{x x} u^{2}+2 \Phi_{x y} u v+\Phi_{y y} v^{2}
$$

should be positive ${ }^{\mathrm{a}}$ for all $u$, $v$ and all $(x, y)$ of $D$.
If $Q$ is strictly positive ${ }^{\mathrm{b}}$, then (3.12.3) holds with inequality, unless all the $x$ and all the $y$ are equal.
(1) The condition is necessary. If $(x, y)$ is in $D$, then $\chi(t)$, defined by (3.12.2), is convex in a neighbourhood of $t=0$. Hence, by Theorem 94, $\chi^{\prime \prime}(0) \geqq 0$, i.e. $Q \geqq 0$.
(2) The condition is sufficient. If

$$
\Sigma q=1, \quad X=\Sigma q x, \quad Y=\Sigma q y
$$

then

$$
\begin{aligned}
& \Phi\left(x_{\nu}, y_{\nu}\right)=\Phi(X, Y)+\left(x_{\nu}-X\right) \Phi_{x}{ }^{0}+\left(y_{\nu}-Y\right) \Phi_{y}{ }^{0} \\
& \quad+\frac{1}{2}\left\{\left(x_{\nu}-X\right)^{2} \Phi_{x x}{ }^{1}+2\left(x_{\nu}-X\right)\left(y_{\nu}-Y\right) \Phi_{x y}{ }^{1}+\left(y_{\nu}-Y\right)^{2} \Phi_{v y}{ }^{1}\right\}
\end{aligned}
$$ where the index 0 indicates the point $(X, Y)$ and the index 1 some point on the line joining this point to $\left(x_{\nu}, y_{\nu}\right)$. It follows that

$$
\Sigma q \Phi(x, y) \geqq \Phi(X, Y)=\Phi(\Sigma q x, \Sigma q y)
$$

If $Q$ is strictly positive, and there is equality, then $x_{\nu}=X$, $y_{\nu}=Y$ for all $\nu$.

$$
\text { a } Q \geqq 0 . \quad \text { b } Q>0 \text { except for } u=v=0
$$

We notice that $Q$ is positive if and only if

$$
\begin{equation*}
\Phi_{x x} \geqq 0, \quad \Phi_{y y} \geqq 0, \quad \Phi_{x x} \Phi_{y y}-\Phi_{x y}{ }^{2} \geqq 0, \tag{3.12.4}
\end{equation*}
$$

and strictly positive if and only if

$$
\begin{equation*}
\Phi_{x x}>0, \quad \Phi_{x x} \Phi_{y y}-\Phi_{x y}^{2}>0 . \tag{3.12.5}
\end{equation*}
$$

It is negative if (3.12.4) holds with the signs of the first two inequalities reversed.

The extension of the definitions and theorems of this section to functions of more than two variables may be left to the reader.
3.13. Generalisations of Hölder's inequality. We may write Hölder's inequality in the form

$$
\begin{equation*}
\mathfrak{A}(a b) \leqq \mathfrak{M}_{F}(a) \mathfrak{M}_{G}(b) \tag{3.13.1}
\end{equation*}
$$

or

$$
\begin{equation*}
\Sigma q a b \leqq F^{-1}\{\Sigma q F(a)\} G^{-1}\{\Sigma q G(b)\} \tag{3.13.2}
\end{equation*}
$$

where $F(x)=x^{r}(r>1)$ and $G(x)=x^{r^{\prime}}, r^{\prime}$ being as usual the index conjugate to $r$ in the sense of $\S 2.8$. If we write

$$
\phi=F^{-1}, \quad \psi=G^{-1}, \quad F(a)=x, \quad G(b)=y, \quad a=\phi(x), \quad b=\psi(y),
$$ we obtain

$$
\begin{equation*}
\Sigma q \phi(x) \psi(y) \leqq \phi(\Sigma q x) \psi(\Sigma q y) . \tag{3.13.3}
\end{equation*}
$$

The simplest case of this is

$$
\frac{1}{2}\left\{\phi\left(x_{1}\right) \psi\left(y_{1}\right)+\phi\left(x_{2}\right) \psi\left(y_{2}\right)\right\} \leqq \phi\left\{\frac{1}{2}\left(x_{1}+x_{2}\right)\right\} \psi\left\{\frac{1}{2}\left(y_{1}+y_{2}\right)\right\},
$$

which expresses the fact that $\phi(x) \psi(y)$ is a concave function of $x$ and $y$. When, as here, $\phi$ and $\psi$ are continuous, it is equivalent to the more general inequality (3.13.3). Hence, reversing the argument (with general $\phi$ and $\psi$ ), we obtain
100. If $F$ and $G$ are continuous and strictly monotonic, then a necessary and sufficient condition that $\mathfrak{A}(a b)$ should be comparable with $\mathfrak{M}_{F}(a) \mathfrak{M}_{G}(b)$ is that $F^{-1}(x) G^{-1}(y)$ should be a concave or convex function of the two variables $x$ and $y$; in the first case (3.13.1) is true, in the second the reverse inequality.

As an example we may take $F(x)=x^{r}, G(y)=y^{s}$. It then follows from Theorems 100 and 99 that

$$
\mathfrak{A}(a b) \leqq \mathfrak{M}_{r}(a) \mathfrak{M}_{s}(b)
$$

if $r>1, s>1$ and $(r-1)(s-1) \geqq 1$. If $r<1, s<1,(1-r)(1-s) \geqq 1$,
the inequality is reversed. These are the only cases of comparability ${ }^{\text {a }}$. The argument excludes the cases $r=0$ and $s=0$, but they may be included by using an exponential instead of a power.

We might look for a more straightforward generalisation of Hölder's inequality. Hölder's inequality asserts that (3.13.1) holds if $f(x)$ and $g(x)$ are inverse positive powers of $x$, and either

$$
\begin{equation*}
F(x)=x f(x), \quad G(x)=x g(x) \tag{a}
\end{equation*}
$$

or

$$
\begin{equation*}
F(x)=\int_{0}^{x} f(t) d t, \quad G(x)=\int_{0}^{x} g(t) d t \tag{b}
\end{equation*}
$$

and we might expect that it would hold for other pairs of inverse $f$ and $g$. The theorem which follows shows that no such extension is possible.
101. Suppose that $f(x)$ is a continuous and strictly increasing function which vanishes for $x=0$ and has a second derivative continuous for $x>0$, and that $g(x)$ is the inverse function (which has necessarily the same properties). Suppose further that $F(x)$ and $G(x)$ are defined either by $(a)$ or by (b), and that (3.13.1) is true for all positive $a, b$. Then $f$ is a power of $x$ and (3.13.1) is Hölder's inequality.

We consider case ( $a$ ). ${ }^{\text {b }}$ If, as in the proof of Theorem 100, we write $\phi$ and $\psi$ for $F^{-1}$ and $G^{-1}$, then $\phi(x) \psi(y)$ must be a concave function of $x$ and $y$. Itfollows from Theorem 99 and (3.12.4)c that $\phi^{\prime \prime} \leqq 0, \psi^{\prime \prime} \leqq 0$ and
(3.13.4) $\quad\left\{\phi^{\prime}(x) \psi^{\prime}(y)\right\}^{2} \leqq \phi(x) \psi(y) \phi^{\prime \prime}(x) \psi^{\prime \prime}(y)$
for all positive $x$ and $y$.
If $\phi(x)=u, \psi(x)=v$, we have

$$
\begin{gathered}
x=F(u)=u f(u), \quad \frac{x}{u}=f(u), \quad u=g\left(\frac{x}{u}\right), \\
x=\frac{x}{u} g\left(\frac{x}{u}\right)=G\left(\frac{x}{u}\right), \quad \frac{x}{u}=\psi(x)=v,
\end{gathered}
$$

and so

$$
\begin{equation*}
\phi(x) \psi(x)=x \tag{3.13.5}
\end{equation*}
$$

${ }^{\text {a }}$ Compare Theorem $44 . \quad{ }^{\text {b }}$ See Cooper (4) for case (b).
c With the appropriate changes of sign.

Hence

$$
\phi^{\prime \prime}(x) \psi(x)+2 \phi^{\prime}(x) \psi^{\prime}(x)+\phi(x) \psi^{\prime \prime}(x)=0,
$$ and so, by (3.13.4),

$$
\left(\phi^{\prime \prime} \psi+\phi \psi^{\prime \prime}\right)^{2}=4 \phi^{\prime 2} \psi^{\prime 2} \leqq 4 \phi \psi \phi^{\prime \prime} \psi^{\prime \prime},
$$

all the arguments being $x$. This is only possible if

$$
\phi^{\prime \prime} \psi=\phi \psi^{\prime \prime}=-\phi^{\prime} \psi^{\prime}, \quad \phi^{\prime \prime} \psi+\phi^{\prime} \psi^{\prime}=0,
$$

or $\phi^{\prime} \psi$ is constant. Hence, by (3.13.5), $x \phi^{\prime} / \phi$ is constant, in which case $\phi$ and the other functions are powers of $x$.
3.14. Some theorems concerning monotonic functions. We collect here some simple theorems which will be useful later. The first characterises monotonic functions as (3.4.2) characterises continuous convex functions.
102. ${ }^{\text {a }}$ A necessary and sufficient condition that
(3.14.1)

$$
(\Sigma p) \phi(\Sigma x) \leqq \Sigma p \phi(x),
$$

for all positive $x$ and $p$, is that $\phi(x)$ should decrease (in the wide sense) for $x>0$. The opposite inequality is similarly characteristic of increasing functions.

There is strict inequality if $\phi(x)$ decreases strictly and there is more than one $x$.
(i) If $\phi$ decreases, $\phi(\Sigma x) \leqq \phi(x)$, whence (3.14.1) follows.
(ii) If in (3.14.1) we take $n=2, x_{1}=x, x_{2}=h, p_{1}=1, p_{2}=p$, we obtain

$$
(1+p) \phi(x+h) \leqq \phi(x)+p \phi(h) .
$$

Making $p \rightarrow 0$, we see that $\phi(x+h) \leqq \phi(x)$.
The case $\phi(x)=x^{\alpha-1}(0<\alpha<1), p=x$, gives Theorem 19.
103. A sufficient condition that

$$
\begin{equation*}
f(\Sigma x) \leqq \Sigma f(x), \tag{3.14.2}
\end{equation*}
$$

for all positive $x$, is that $x^{-1} f(x)$ should decrease. There is strict inequality if $x^{-1} f(x)$ decreases strictly and there is more than one $x$.

For if we write $f(x)=x \phi(x)$, then (3.14.2) becomes (3.14.1) with $p=x$. The condition is not necessary, since (3.14.2) is satisfied by any $f(x)$ for which

$$
f(x)>0, \quad \operatorname{Max} f(x) \leqq 2 \operatorname{Min} f(x) ;
$$

for example $f(x)=3+\cos x$.

[^23]104. If
(3.14.3)
$$
\phi(\Sigma p x) \leqq \Sigma p \phi(x)
$$
for all positive $x$ and $p$, then $\phi(x)$ is a multiple of $x$.
If we take $n=2, x_{1}=x, x_{2}=y, p_{1}=y / 2 x$ and $p_{2}=\frac{1}{2}$ in (3.14.3), we obtain
$$
\frac{\phi(y)}{y} \leqq \frac{\phi(x)}{x} .
$$

Since we can interchange $x$ and $y, \phi / x$ is constant.

### 3.15. Sums with an arbitrary function: generalisations

 of Jensen's inequality. We may define 'sums' involving an arbitrary function $\phi$ as well as means. We write$$
\mathbb{S}_{\phi}(a)=\phi^{-1}\{\Sigma \phi(a)\} .
$$

Here $\phi(x)$ is continuous and strictly monotonic, as in § 3.1 ; but it is necessary now to assume rather more, since $\Sigma \phi(a)$ is not, like $\Sigma q \phi(a)$, a mean of the values of $\phi(a)$, and so necessarily a possible value of $\phi(x)$. We therefore assume that $\phi(x)$ is positive for all positive $x$ and tends to $\infty$ either when $x \rightarrow 0$ or when $x \rightarrow \infty$. We shall also assume that the $a$ are all positive, leaving the reader to make the modifications appropriate when any $a$ is zero ${ }^{2}$.
105. ${ }^{\text {b }}$ If $\psi$ and $\chi$ are continuous, positive, and strictly monotonic, then $\Im_{\psi}$ and $\varsigma_{\chi}$ are comparable whenever (1) $\psi$ and $\chi$ vary in opposite directions, or (2) $\psi$ and $\chi$ vary in the same direction and $\chi / \psi$ is monotonic.

In case (1),

$$
\begin{equation*}
\mathfrak{S}_{\psi} \leqq \mathscr{S}_{x} \tag{3.15.1}
\end{equation*}
$$

if $\psi$ decreases and $\chi$ increases. In case (2), (3.15.1) is true if $\chi / \psi$ decreases. There is equality in case (1), and, if $\chi / \psi$ is strictly monotonic, also in case (2), only when there is only one $a$.
a Suppose, for example, that $\phi(x)=x^{r}$, where $r>0$ (the case of $\S 2.10$ ). Then $\phi(0)=0$, and we need make no distinction between two such systems of the $a$ as $(1,1)$ and $(1,1,0)$. If $\phi(0)$ were positive it would be necessary to distinguish, and the discrimination of the cases of equality in Theorem 105 would become tedious. If $\phi(x) \rightarrow \infty$ when $x \rightarrow 0$, then $\mathfrak{S}_{\phi}(a)=0$ whenever any $a$ is zero.
${ }^{b}$ The substance of this theorem is due to Cooper (2).

In case (1), when $\chi$ increases,

$$
\mathfrak{S}_{\chi}(a) \geqq \chi^{-1}\{\chi(\operatorname{Max} a)\}=\operatorname{Max} a
$$

and similarly $\mathfrak{S}_{\psi}(a) \leqq \operatorname{Min} a$.
In case (2), suppose that $\psi$ and $\chi$ increase, and write

$$
\psi(a)=x, \quad a=\psi^{-1}(x), \quad \chi \psi^{-1}=f .
$$

Then (3.15.1) reduces to (3.14.2), and is true if $x^{-1} f(x)$ decreases, i.e. if

$$
\frac{f\{\psi(x)\}}{\psi(x)}=\frac{\chi(x)}{\psi(x)}
$$

decreases. If $\psi$ and $\chi$ decrease, (3.15.1) reduces to (3.14.2) reversed, and is true if $f / x$ increases, or if $\chi / \psi$ decreases.

The reader will have no difficulty in distinguishing the cases of equality. The case $\psi=x^{s}, \chi=x^{r}$ gives Theorem 19.

We may also define weighted sums analogous to those of § 2.10 (iv), viz.

$$
\mathfrak{I}_{\phi}(a)=\phi^{-1}\{\Sigma p \phi(a)\},
$$

where the $p$ are arbitrary positive numbers. $\mathfrak{T}_{\phi}$ reduces to $\mathfrak{M}_{\phi}$ if $\Sigma p=1$, to $\mathfrak{S}_{\phi}$ if every $p$ is 1 .
3.16. Generalisations of Minkowski's inequality. If $\phi(x)=x^{r}$, where $r>1$, we have

$$
\begin{align*}
\mathfrak{M}_{\phi}\left(\frac{a+b}{2}\right) & \leqq \frac{1}{2}\left\{\mathfrak{M}_{\phi}(a)+\mathfrak{M}_{\phi}(b)\right\},  \tag{3.16.1}\\
\mathfrak{S}_{\phi}\left(\frac{a+b}{2}\right) & \leqq \frac{1}{2}\left\{\Im_{\phi}(a)+\mathfrak{S}_{\phi}(b)\right\} \\
\mathfrak{I}_{\phi}\left(\frac{a+b}{2}\right) & \leqq \frac{1}{2}\left\{\mathfrak{V}_{\phi}(a)+\mathfrak{I}_{\phi}(b)\right\},
\end{align*}
$$

all these inequalities being essentially equivalent and included in Theorem 24.

The inequalities are not equivalent for general $\phi$; all of them are of the form

$$
\begin{equation*}
2 \phi^{-1}\left\{\Sigma p \phi\left(\frac{a+b}{2}\right)\right\} \leqq \phi^{-1}\{\Sigma p \phi(a)\}+\phi^{-1}\{\Sigma p \phi(b)\} \tag{3.16.4}
\end{equation*}
$$

but the differences between the weights $p$ are now significant. In (3.16.1), $\Sigma p=1$; in (3.16.2), $p=1$; in (3.16.3), the $p$ are any
positive numbers. We call these three cases the cases (I), (II), and (III). In discussing them, we shall suppose that
for $x>0$.

$$
\phi>0, \quad \phi^{\prime}>0
$$

The inequality (3.16.4) asserts that $\phi^{-1}\{\Sigma p \phi(x)\}$ is, for given $p$, a convex function of the $n$ variables $x_{1}, x_{2}, \ldots, x_{n}$; or, after §3.12, ${ }^{\text {a }}$ that

$$
\chi(t)=\phi^{-1}\{\Sigma p \phi(x+u t)\}
$$

where the $x, p, u$ are fixed and the $x$ and $p$ positive, is convex in $t$ for all $t$ for which all $x+u t$ are positive. If $\phi$ is twice differentiable this condition is, by Theorem 94, equivalent to $\chi^{\prime \prime}(0) \geqq 0$. A straightforward calculation shows that

$$
\begin{equation*}
\left\{\phi^{\prime}(\chi)\right\}^{3} \chi^{\prime \prime}=\left\{\phi^{\prime}(\chi)\right\}^{2} \Sigma p u^{2} \phi^{\prime \prime}(x)-\phi^{\prime \prime}(\chi)\left\{\Sigma p u \phi^{\prime}(x)\right\}^{2}, \tag{3.16.5}
\end{equation*}
$$

where

$$
\begin{equation*}
\chi=\chi(0)=\phi^{-1}\{\Sigma p \phi(x)\} \tag{3.16.6}
\end{equation*}
$$

and $\chi^{\prime \prime}=\chi^{\prime \prime}(0)$. We have therefore to consider in what circumstances

$$
\begin{equation*}
\left\{\phi^{\prime}(\chi)\right\}^{2} \Sigma p u^{2} \phi^{\prime \prime}(x)-\phi^{\prime \prime}(\chi)\left\{\Sigma p u \phi^{\prime}(x)\right\}^{2} \geqq 0 \tag{3.16.7}
\end{equation*}
$$

It is easy to see that (3.16.7) cannot be true generally without restriction on the sign of $\phi^{\prime \prime}$. Suppose for example that $\phi^{\prime}>0$ and that $\phi^{\prime \prime}$ is continuous and sometimes negative. We can then choose $x_{1}$ and $x_{2}$ so that $\phi^{\prime \prime}\left(x_{1}\right)<0, \phi^{\prime \prime}\left(x_{2}\right)<0$, and $u_{1}$ and $u_{2}$ so that

$$
p_{1} u_{1} \phi^{\prime}\left(x_{1}\right)+p_{2} u_{2} \phi^{\prime}\left(x_{2}\right)=0
$$

In this case (3.16.7), for $n=2$, reduces to

$$
\left\{\phi^{\prime}(\chi)\right\}^{2}\left\{p_{1} u_{1}^{2} \phi^{\prime \prime}\left(x_{1}\right)+p_{2} u_{2}^{2} \phi^{\prime \prime}\left(x_{2}\right)\right\} \geqq 0,
$$

which is false. We shall therefore suppose in what follows that

$$
\phi>0, \phi^{\prime}>0, \phi^{\prime \prime}>0
$$

We can write (3.16.7) in the form

$$
\begin{equation*}
\psi(\chi)=\frac{\phi^{\prime 2}(\chi)}{\phi^{\prime \prime}(\chi)} \geqq \frac{\left\{\Sigma p u \phi^{\prime}(x)\right\}^{2}}{\Sigma p u^{2} \phi^{\prime \prime}(x)} \tag{3.16.8}
\end{equation*}
$$

[^24]Now, by Theorem 7,
(3.16.9)

$$
\left(\Sigma p u \phi^{\prime}\right)^{2}=\left\{\Sigma \sqrt{ }\left(p \phi^{\prime \prime}\right) u \cdot \sqrt{ }\left(\frac{p \phi^{\prime 2}}{\phi^{\prime \prime}}\right)\right\}^{2} \leqq \Sigma p u^{2} \phi^{\prime \prime} \Sigma p \frac{\phi^{\prime 2}}{\phi^{\prime \prime}} .
$$

Hence (3.16.8) is certainly true for all $x, u$ if

$$
\begin{equation*}
\psi(x) \geqq \Sigma p \frac{\left\{\phi^{\prime}(x)\right\}^{2}}{\phi^{\prime \prime}(x)}=\Sigma p \psi(x), \tag{3.16.10}
\end{equation*}
$$

for all $x$. Further, there is equality in (3.16.9) if

$$
u=\frac{\phi^{\prime}(x)}{\phi^{\prime \prime}(x)} \quad(\nu=1,2, \ldots, n),
$$

so that (3.16.10) is both a sufficient and a necessary condition for the truth of (3.16.8). Finally, if we write $y=\phi(x)$ and

$$
\begin{equation*}
\Phi(y)=\psi(x)=\psi\left\{\phi^{-1}(y)\right\}=\frac{\left\{\phi^{\prime}\left[\phi^{-1}(y)\right]\right\}^{2}}{\phi^{\prime \prime}\left\{\phi^{-1}(y)\right\}}, \tag{3.16.11}
\end{equation*}
$$

(3.16.10) assumes the form

$$
\begin{equation*}
\Phi(\Sigma p y) \geqq \Sigma p \Phi(y) . \tag{3.16.12}
\end{equation*}
$$

We now consider the three cases (I), (II), (III) separately.
(i) In case (I), (3.16.12) is true if and only if $\Phi(y)$ is a concave function of $y$.
(ii) In case (II), (3.16.12) is (3.14.2), reversed and with $y, \Phi$ for $x, f$. A sufficient (though not a necessary) condition is that $\Phi / y$ is an increasing function of $y$ or, what is the same thing, that $\phi \phi^{\prime \prime} \mid \phi^{\prime 2}$ is a decreasing function of $x$.
(iii) In case (III), (3.16.12) is (3.14.3), with the appropriate variations. It can be true generally only when $\Phi(y)$ is a multiple of $y$, or when $\phi \phi^{\prime \prime} \mid \phi^{\prime 2}$ is constant, in which case $\phi$ is of one of the forms

$$
\begin{equation*}
(a x+b)^{c}\left(a>0, c>1^{\mathrm{a}}\right), \quad e^{a x+b} \tag{3.16.13}
\end{equation*}
$$

In these cases it is true.
There are alternative forms of the conditions (i) and (ii) which show better their relations to one another. We shall suppose $\phi^{\prime \prime \prime \prime}$ continuous, as we may do without affecting seriously the interest
of the results. Then from

$$
\Phi(y)=\frac{\phi^{\prime 2}(x)}{\phi^{\prime \prime}(x)}
$$

we deduce

$$
\begin{gathered}
\Phi^{\prime}(y)=\frac{d}{d x} \frac{\phi^{\prime}(x)}{\phi^{\prime \prime}(x)}+1, \\
\Phi^{\prime \prime}(y)=\frac{1}{\phi^{\prime}(x)} \frac{d^{2}}{d x^{2}} \frac{\phi^{\prime}(x)}{\phi^{\prime \prime}(x)} .
\end{gathered}
$$

Hence $\Phi(y)$ is concave if and only if $\phi^{\prime}(x) / \phi^{\prime \prime}(x)$ is concave, or, what is the same thing, if $\phi^{\prime} \phi^{\prime \prime \prime} / \phi^{\prime 2}$ is increasing. These are alternative forms of (i), and an alternative form of (ii) is ' $\phi / \phi^{\prime}$ is convex'.

Summing up our conclusions, we have
106. ${ }^{\text {a }}$ Suppose that $\phi^{\prime \prime \prime \prime}$ is continuous and that $\phi>0, \phi^{\prime}>0$, $\phi^{\prime \prime}>0$. Then
(i) it is necessary and sufficient for (3.16.1) that $\phi^{\prime} / \phi^{\prime \prime}$ should be concave, or $\phi^{\prime} \phi^{\prime \prime \prime} / \phi^{\prime 2}$ increasing;
(ii) it is sufficient (but not necessary) for (3.16.2) that $\phi / \phi^{\prime}$ should be convex, or $\phi \phi^{\prime \prime} / \phi^{\prime 2}$ decreasing;
(iii) it is necessary and sufficient for (3.16.3) that $\phi$ should be one of the functions (3.16.13).

We leave it for the reader to formulate the variations of this theorem, when (for example) $\phi>0, \phi^{\prime}<0, \phi^{\prime \prime}>0$, or when the inequalities are reversed. It is instructive to verify that (i) is satisfied (from a certain $x$ onwards) when $\phi=x^{p} / \log x$, where $p>1$, but not when $\phi=x^{p} \log x$, while for (ii) the situation is reversed.
3.17. Comparison of sets. Theorem 105 asserts that

$$
\mathbb{S}_{\psi}(a) \leqq \mathbb{S}_{\chi}(a)
$$

for a given pair of functions $\psi$ and $\chi$ and all $a$. The theorems of this section are of a different type, involving given sets (a) and ( $a^{\prime}$ ) and a variable function $\phi$. We consider the conditions under which

$$
\mathbb{S}_{\phi}\left(a^{\prime}\right) \leqq \mathbb{S}_{\phi}(a),
$$

[^25]or, what is the same thing for increasing $\phi$,
(3.17.1)
$$
\phi\left(a_{1}{ }^{\prime}\right)+\phi\left(a_{2}{ }^{\prime}\right)+\ldots+\phi\left(a_{n}{ }^{\prime}\right) \leqq \phi\left(a_{1}\right)+\phi\left(a_{2}\right)+\ldots+\phi\left(a_{n}\right)
$$
for given $a$ and $a^{\prime}$ and all $\phi$ of a certain class.
107. Suppose that the sets ( $a$ ) and ( $a^{\prime}$ ) are arranged in descending order of magnitude. Then a necessary and sufficient condition that (3.17.1) should be true for all continuous and increasing $\phi$ is that
$$
a_{\nu}^{\prime} \leqq a_{\nu} \quad(\nu=1,2, \ldots, n) .
$$

The sufficiency of the condition is obvious. To prove it necessary, suppose that $a_{\mu}{ }^{\prime}>a_{\mu}$ for some $\mu$, that $a_{\mu}<b<a_{\mu}{ }^{\prime}$, and that $\phi^{*}(x)$ is defined by

$$
\begin{gathered}
\phi^{*}(x)=0(x \leqq b), \quad \phi^{*}(x)=1(x>b) . \\
\Sigma \phi^{*}\left(a^{\prime}\right) \geqq \mu>\mu-1 \geqq \Sigma \phi^{*}(a) .
\end{gathered}
$$

Then
Hence (3.17.1) is false for $\phi^{*}$, and therefore also false for an appropriately chosen continuous increasing approximation to $\phi^{*}$.

Ournext theorem is connected with the theorems of §§ 2.18-2.20.
108. In order that (3.17.1) should be true for all continuous convex $\phi$, it is necessary and sufficient either that (1) $\left(a^{\prime}\right) \prec(a)$, i.e. that ( $a^{\prime}$ ) is majorised by ( $a$ ) in the sense of § 2.18, or that (2) ( $a^{\prime}$ ) is an average of ( $a$ ) in the sense of §2.20.

If these conditions are satisfied, and $\phi^{\prime \prime}(x)$ exists for all $x$, and is positive, then equality can occur in (3.17.1) only when the sets (a) and ( $a^{\prime}$ ) are identicala.

We have proved already (Theorem 46) that the two conditions are equivalent. It is therefore enough to prove that the first is necessary and the second sufficient. We may suppose (a) and ( $a^{\prime}$ ) arranged in descending order.
(i) Condition (1) is necessary. Condition (1) asserts that
(3.17.2) $\quad a_{1}{ }^{\prime}+a_{2}{ }^{\prime}+\ldots+a_{\nu}{ }^{\prime} \leqq a_{1}+a_{2}+\ldots+a_{\nu} \quad(\nu=1,2, \ldots, n)$, with equality in the case $\nu=n$.
The functions $x$ and $-x$ are both continuous and convex in any

[^26]interval. Hence, if (3.17.1) is true, $\Sigma a^{\prime} \leqq \Sigma a$ and $\Sigma\left(-a^{\prime}\right) \leqq \Sigma(-a)$, i.e. $\Sigma a^{\prime}=\Sigma a$, which is (3.17.2) for $\nu=n$, with equality.

Next, let

$$
\phi(x)=0 \quad\left(x \leqq a_{\nu}\right), \quad \phi(x)=x-a_{\nu}\left(x>a_{\nu}\right) .
$$

Then $\phi(x)$ is continuous and convex in any interval, and $\phi(x) \geqq 0$, $\phi(x) \geqq x-a_{\nu}$. Hence
$a_{1}{ }^{\prime}+a_{2}{ }^{\prime}+\ldots+a_{\nu}{ }^{\prime}-\nu a_{\nu} \leqq \Sigma \phi\left(a^{\prime}\right) \leqq \Sigma \phi(\alpha)=a_{1}+a_{2}+\ldots+a_{\nu}-\nu a_{\nu}$, which is (3.17.2).
(ii) Condition (2) is sufficient. If ( $a^{\prime}$ ) is an average of $(a)$, we have

$$
a_{\mu}^{\prime}=p_{\mu 1} a_{1}+p_{\mu 2} a_{2}+\ldots+p_{\mu n} a_{n}
$$

where

$$
p_{\mu \nu} \geqq 0, \quad \sum_{\mu=1}^{n} p_{\mu \nu}=1, \quad \sum_{\nu=1}^{n} p_{\mu \nu}=1
$$

for all $\mu$ and $\nu$. If $\phi$ is convex, then

$$
\begin{equation*}
\phi\left(a_{\mu}^{\prime}\right) \leqq p_{\mu 1} \phi\left(a_{1}\right)+\ldots+p_{\mu n} \phi\left(a_{n}\right) \tag{3.17.3}
\end{equation*}
$$

and (3.17.1) follows by summation.
(iii) If there is equality in (3.17.1), there must be equality in each of (3.17.3).

If $\phi^{\prime \prime}(x)>0$, and every $p_{\mu \nu}$ is positive, then it follows from Theorem 95 that all the $a$ are equal, in which case all the $a^{\prime}$ are also equal and the common values are the same.

In general, however, some of the $p_{\mu \nu}$ will be zero. We shall say that $a_{\mu}{ }^{\prime}$ and $a_{\nu}$ are immediately connected if $p_{\mu \nu}>0$, i.e. if $a_{\nu}$ occurs effectively in the formula for $a_{\mu}{ }^{\prime}$; and that any two elements (whether $a$ or $a^{\prime}$ ) are connected if they can be joined by a chain of elements in which each consecutive pair is immediately connected.

Consider now the complete set $C$ of elements connected with $a_{1}$. We may write this set (changing the numeration of the elements if necessary) as

$$
\begin{equation*}
a_{1}, a_{2}, \ldots, a_{r}, \quad a_{1}^{\prime}, a_{2}^{\prime}, \ldots, a_{s}^{\prime} \tag{C}
\end{equation*}
$$

the $a^{\prime}$ of $C$ involve the $a$ of $C$, and no other $a$, and no other $a^{\prime}$ involves an $a$ of $C$. Hence, using the sum-properties of the $p$,

$$
s=\sum_{\mu=1}^{s} \sum_{\nu=1}^{r} p_{\mu \nu}=\sum_{\nu=1}^{r} \sum_{\mu=1}^{s} p_{\mu \nu}=r
$$

so that $C$ contains just as many $a^{\prime}$ as $a$. It follows from Theorem 95, and from the equality in (3.17.3), that all $a$ immediately connected with an $a^{\prime}$ are equal to that $a^{\prime}$. Hence all connected $a$ and $a^{\prime}$ are equal, and $C$ contains $r$ elements of each set, all equal to $a_{1}$.

We now repeat the argument, starting from $a_{r+1}$, and we conclude that both ( $a$ ) and ( $a^{\prime}$ ) consist of a certain number of blocks of equal elements, the values of the elements in corresponding blocks being the same.

Incidentally we have proved
109. If $\phi^{\prime \prime}(x)>0, p_{\mu \nu}>0, \sum_{\mu} p_{\mu \nu}=1, \sum_{\nu} p_{\mu \nu}=1$, and $a_{\mu}^{\prime}=\sum_{\nu} p_{\mu \nu} a_{\nu}$, then

$$
\begin{equation*}
\Sigma \phi\left(a^{\prime}\right)<\Sigma \phi(a) \tag{3.17.4}
\end{equation*}
$$

unless all the $a$ and $a^{\prime}$ are equal.
If all the $a^{\prime}$ are equal, (3.17.4) is a special case of Theorem 95.
A special case of Theorem 108 which is often useful is
110. If $\phi(x)$ is continuous and convex, and $\left|h^{\prime}\right| \leqq|h|$, then

$$
\begin{equation*}
\phi\left(x-h^{\prime}\right)+\phi\left(x+h^{\prime}\right) \leqq \phi(x-h)+\phi(x+h) . \tag{3.17.5}
\end{equation*}
$$

3.18. Furthergeneral properties of convexfunctions. We have assumed since $\S 3.6$ that $\phi(x)$ is continuous. We shall now discard this hypothesis and consider the direct consequences of (3.5.1). The general lesson of the theorems which follow will be that a convex function is either very regular or very irregular, and in particular that a convex function which is not 'entirely irregular' is necessarily continuous (so that the hypothesis of continuity is a good deal less restrictive than might have been expected).
111. Suppose that $\phi(x)$ is convex in the open interval $(H, K)$, and bounded above in some interval $i$ interior to $(H, K)$. Then $\phi(x)$ is continuous in the open interval $(H, K)$. Further, $\phi(x)$ has everywhere left-hand and right-hand derivatives; the right-hand derivative is not less than the left-hand derivative; and both derivatives increase with $x$.

It follows that a discontinuous convex function is unbounded in every interval.

We prove first that $\phi(x)$ is bounded above in every interval interior to $(H, K)$. The kernel of the proof is this. The argument of § 3.6 shows that

$$
\phi(\Sigma q x) \leqq \Sigma q \phi(x)
$$

for any rational $q$; it was only in the passage to irrational $q$ that we used the hypothesis of continuity. Suppose now that $i$ is ( $h, k$ ) and that the upper bound of $\phi$ in $i$ is $G$. It is enough to prove $\phi$ bounded above in $(l, h)$ and ( $k, m$ ), where $l$ and $m$ are any numbers such that $H<l<h<k<m<K$. If $x$ is in ( $l, h$ ), we can find a $\xi$ in $i$ so that $x$ divides $(l, \xi)$ rationally, and then $\phi(x)$ must lie below a bound depending on $\phi(l)$ and $G$, and so be bounded above in $(l, h)$. Similarly, it must be bounded above in $(k, m)$.

To state the argument precisely, let $h$ be the left-hand end of $i$ and $G$ the upper bound of $\phi$ in $i$, and suppose that

$$
H<l<x<h
$$

We can choose integers $m$ and $n>m$ so that

$$
\xi=l+\frac{n}{m}(x-l)
$$

lies in $i$, and then

$$
\begin{aligned}
\phi(x)=\phi\left\{\frac{m \xi+(n-m) l}{n}\right\} & \leqq \frac{m}{n} \phi(\xi)+\frac{n-m}{n} \phi(l) \\
& \leqq \frac{m}{n} G+\frac{n-m}{n} \phi(l) \leqq \operatorname{Max}\{G, \phi(l)\} .
\end{aligned}
$$

Hence $\phi(x)$ is bounded above in $(l, h)$.
In proving the remainder of the theorem we may restrict ourselves to an interval ( $H^{\prime}, K^{\prime}$ ) interior to ( $H, K$ ), or, what is the same thing, we may suppose $\phi$ bounded above in the original interval. Suppose then that $\phi(x) \leqq G$ in $(H, K)$, that $H<x<K$, that $m$ and $n>m$ are positive integers, and that $\delta$ is a number (positive or negative) small enough to leave $x+n \delta$ inside ( $H, K$ ). Then

$$
\phi(x+m \delta)=\phi\left\{\frac{m(x+n \delta)+(n-m) x}{n}\right\} \leqq \frac{m}{n} \phi(x+n \delta)+\frac{n-m}{n} \phi(x),
$$

or

$$
\frac{\phi(x+n \delta)-\phi(x)}{n} \geqq \frac{\phi(x+m \delta)-\phi(x)}{m} \quad(m<n) .
$$

Replacing $\delta$ by $-\delta$, and combining the two inequalities, we obtain

$$
\begin{align*}
& \frac{\phi(x+n \delta)-\phi(x)}{n} \geqq \frac{\phi(x+m \delta)-\phi(x)}{m}  \tag{3.18.1}\\
& \quad \geqq \frac{\phi(x)-\phi(x-m \delta)}{m} \geqq \frac{\phi(x)-\phi(x-n \delta)}{n}
\end{align*}
$$

(the central inequality following directly from the convexity of $\phi$ ).
If in (3.18.1) we take $m=1$, and remember that $\phi \leqq G$, we find that

$$
\begin{equation*}
\frac{G-\phi(x)}{n} \geqq \phi(x+\delta)-\phi(x) \geqq \phi(x)-\phi(x-\delta) \geqq \frac{\phi(x)-G}{n} . \tag{3.18.2}
\end{equation*}
$$

We now suppose that $\delta \rightarrow 0$ and $n \rightarrow \infty$, but so that $x \pm n \delta$ remains inside the interval. It then follows from (3.18.2) that $\phi(x+\delta)$ and $\phi(x-\delta)$ tend to $\phi(x)$, and so that $\phi$ is continuous.

We next suppose $\delta>0$, and replace $\delta$ in (3.18.1) by $\delta / n$. We have then

$$
\begin{align*}
\frac{\phi(x+\delta)-\phi(x)}{\delta} & \geqq \frac{\phi\left(x+\delta^{\prime}\right)-\phi(x)}{\delta^{\prime}}  \tag{3.18.3}\\
& \geqq \frac{\phi(x)-\phi\left(x-\delta^{\prime}\right)}{\delta^{\prime}} \geqq \frac{\phi(x)-\phi(x-\delta)}{\delta},
\end{align*}
$$

where $\delta^{\prime}=m \delta / n$ is any rational multiple of $\delta$ less than $\delta$. Since $\phi$ is continuous, (3.18.3) is true for any $\delta^{\prime}<\delta$. It follows that the quotients on the extreme left and right decrease and increase respectively when $\delta$ decreases to zero, and so that each tends to a limit. Hence $\phi$ possesses right-handed and left-handed derivatives $\phi_{r}{ }^{\prime}$ and $\phi_{l}{ }^{\prime}$, and $\phi_{l}{ }^{\prime} \leqq \phi_{r}{ }^{\prime}$.

Finally, we may write $x-\delta^{\prime}=x_{1}, x=x_{2}, x+\delta=x_{3}\left(\right.$ or $x-\delta=x_{1}$, $x=x_{2}, x+\delta^{\prime}=x_{3}$ ), when (3.18.3) gives

$$
\frac{\phi\left(x_{3}\right)-\phi\left(x_{2}\right)}{x_{3}-x_{2}} \geqq \frac{\phi\left(x_{2}\right)-\phi\left(x_{1}\right)}{x_{2}-x_{1}} .
$$

$A$ fortiori, if $x_{1}<x_{2}<x_{3}<x_{4}$, we have

$$
\frac{\phi\left(x_{4}\right)-\phi\left(x_{3}\right)}{x_{4}-x_{3}} \geqq \frac{\phi\left(x_{2}\right)-\phi\left(x_{1}\right)}{x_{2}-x_{1}}
$$

Making $x_{3} \rightarrow x_{4}, x_{2} \rightarrow x_{1}$, we obtain

$$
\begin{equation*}
\phi_{r}^{\prime}\left(x_{4}\right) \geqq \phi_{l}^{\prime}\left(x_{4}\right) \geqq \phi_{r}^{\prime}\left(x_{1}\right) \geqq \phi_{l}^{\prime}\left(x_{1}\right), \tag{3.18.4}
\end{equation*}
$$

which completes the proof of the theorem.
It is plain from what precedes that

$$
\begin{equation*}
\phi_{l}^{\prime}\left(x_{4}\right) \geqq \frac{\phi\left(x_{3}\right)-\phi\left(x_{2}\right)}{x_{3}-x_{2}} \geqq \phi_{r}^{\prime}\left(x_{1}\right) \tag{3.18.5}
\end{equation*}
$$

if $x_{1} \leqq x_{2}<x_{3} \leqq x_{4}$.
Theorem 111 asserts nothing about the existence of an ordinary differential coefficient $\phi^{\prime}(x)$. It is however easy to prove that $\phi^{\prime}(x)$ exists except perhaps for an enumerable set of values of $x$. The function $\phi_{l}^{\prime}(x)$, being monotonic, is continuous except perhaps in such a set. If it is continuous at $x_{1}$ then, by (3.18.4), $\phi_{r}{ }^{\prime}\left(x_{1}\right)$ lies between $\phi_{l}^{\prime}\left(x_{1}\right)$ and $\phi_{l}^{\prime}\left(x_{4}\right)$, which tends to $\phi_{l}^{\prime}\left(x_{1}\right)$ when $x_{4} \rightarrow x_{1}$. Hence $\phi_{r}^{\prime}\left(x_{1}\right)=\phi_{l}^{\prime}\left(x_{1}\right)$, and $\phi^{\prime}(x)$ exists for $x=x_{1}$.

It is also plain from (3.18.5) that, if $\phi(x)$ is continuous and convex in an open interval $(a, b)$, then

$$
\left|\frac{\phi\left(x^{\prime}\right)-\phi(x)}{x^{\prime}-x}\right|
$$

is bounded for all $x$ and $x^{\prime}$ of any closed sub-interval of $(a, b)$.
3.19. Further properties of continuous convexfunctions. We now suppose $\phi(x)$ convex and continuous. It follows from (3.18.5) that if $H<\xi<K$ and

$$
\phi_{l}^{\prime}(\xi) \leqq \lambda \leqq \phi_{r}^{\prime}(\xi)
$$

then the line

$$
\begin{equation*}
y-\phi(\xi)=\lambda(x-\xi) \tag{3.19.1}
\end{equation*}
$$

will lie wholly under (on or below) the curve. In other words
112. If $\phi(x)$ is convex and continuous then there is at least one line through every point of the curve $y=\phi(x)$ which lies wholly under the curve.

A line through a point of a curve which lies wholly on one side of the curve (under or over) is called a Stützgerade of the curve. If $\phi(x)$ is concave, then the graph of $\phi(x)$ has at every point a Stützgerade over the curve. If $\phi(x)$ is both convex and concave the two lines must coincide and $\phi(x)$ must be linear.

It is easy to see the truth of Theorem 112 directly. If $\xi<x<x^{\prime}$, and $P, Q, Q^{\prime}$ are the points on the curve corresponding to $\xi, x, x^{\prime}$, then $P Q$ lies under $P Q^{\prime}$, and the slope of $P Q$ decreases as $x$ approaches $\xi$, and so tends to a limit $\nu$. Similarly, if $x<\xi$ and $x$ tends to $\xi$, the slope of $Q P$ increases to a limit $\mu$. If $\mu$ were greater than $\nu$, and $x_{1}, x_{2}$ were respectively less than and greater than $\xi$, and sufficiently near to $\xi$, then $P$ would lie above $P_{1} P_{2}$, in contradiction to the convexity of the curve. Hence $\mu \leqq \nu$, and (3.19.1) lies under the curve if $\lambda$ has any value between $\mu$ and $\nu$ inclusive.

In this proof we do not appeal to Theorem 111, but the proof depends on just those geometrical ideas which underlie the more formal and analytical argument of § 3.18.

Suppose now, conversely, that $\phi(x)$ is continuous and has the property assertedin Theorem 112. If $x_{1}$ and $x_{2}$ are two values of $x$, $P_{1}$ and $P_{2}$ the corresponding points on the curve, and $P$ the point corresponding to $\xi=\frac{1}{2}\left(x_{1}+x_{2}\right)$, then both $P_{1}$ and $P_{2}$ lie over a certain line through $P$, and the middle point of $P_{1} P_{2}$ lies over $P$. Hence $\phi(x)$ is convex.

We have thus proved that the property of Theorem 112 affords a necessary and sufficient condition for the convexity of a continuous function, and might be used as an alternative definition of convexity. That is to say, we might define convexity, for continuous functions, as follows: a continuous function $\phi(x)$ is said to be convex in $(H, K)$ if to any $\xi$ of $(H, K)$ corresponds a number $\lambda=\lambda(\xi)$ such that

$$
\phi(\xi)+\lambda(x-\xi) \leqq \phi(x)
$$

for all $x$ of $(H, K)$.
This definition of a convex function is quite as 'natural' as that implied in (3.5.2), and it is interesting to deduce some of the characteristic properties of continuous convex functions directly from it. For example, the inequality (3.4.2) may be proved as follows ${ }^{\text {a }}$.

Writing as usual

$$
\mathfrak{A}(b)=\Sigma q b,
$$

and taking $\xi=\mathfrak{A}(a)$, a value which lies in the interval of variation of the $a$, we have

$$
\begin{gathered}
\phi\{\mathfrak{A}(a)\}+\lambda(a-\xi) \leqq \phi(a) \\
\text { a Jessen (2). }
\end{gathered}
$$

for a certain $\lambda=\lambda(\xi)$ and all $a$. Performing the operation $\mathfrak{A}$ on each side, we obtain
or

$$
\begin{gathered}
\phi\{\mathfrak{A}(a)\}+\lambda\{\mathfrak{A}(a)-\xi\} \leqq \mathfrak{A}\{\phi(a)\}, \\
\phi\{\mathfrak{A}(a)\} \leqq \mathfrak{A}\{\phi(a)\},
\end{gathered}
$$

which is (3.4.2). It is instructive to compare this argument with that of § 3.10 (ii).
3.20. Discontinuous convex functions. Discontinuous convex functions are, by Theorem 111, unbounded in every interval, and their existence has not been proved except under the assumption of Zermelo's Axiom or (what is for our purpose equivalent) the assumption that the continuum can be well-ordered.

If

$$
\begin{gather*}
f(x+y)=f(x)+f(y)  \tag{3.20.1}\\
f(2 x)=2 f(x) \\
2 f\left(\frac{x+y}{2}\right)=f(x+y)=f(x)+f(y)
\end{gather*}
$$

then
and
Thus a solution of (3.20.1) is certainly convex.
It was proved by Hamel (1) that, if Zermelo's Axiom is true, there exist bases $[\alpha, \beta, \gamma, \ldots]$ for the real numbers, that is to say, sets of real numbers $\alpha, \beta, \gamma, \ldots$ such that every real $a$ is expressible uniquely in the form of a finite sum

$$
x=a \alpha+b \beta+\ldots+l \lambda
$$

with rational coefficients $a, b, \ldots, l$. If we assume this, we can at once write down discontinuous solutions of (3.20.1); we give $f(x)$ arbitrary values $f(\alpha), f(\beta), \ldots$ for $x=\alpha, \beta, \ldots$, and define $f(x)$ generally by

$$
f(x)=a f(\alpha)+b f(\beta)+\ldots+l f(\lambda)
$$

Then, if $y=a^{\prime} \alpha+\ldots$, we have

$$
f(x+y)=f\left\{\left(a+a^{\prime}\right) \alpha+\ldots\right\}=\left(a+a^{\prime}\right) f(\alpha)+\ldots=f(x)+f(y)
$$

For more detailed study of the properties of convex functions, of the solutions of the equation (3.20.1) and of inequalities associated with it, we may refer to Darboux (1), Fréchet (1, 2), F. Bernstein (1), Bernstein and Doetsch (1), Blumberg (1), Sierpiński (1, 2), Cooper (3), and Ostrowski (1). Blumberg and Sierpiński prove that any convex measurable function is continuous, and Ostrowski obtains a still more general result.

[^27]
## MISCELLANEOUS THEOREMS AND EXAMPLES

113. If $\alpha$ is a constant, $\alpha \neq 0$, and $\chi=\alpha \phi$, then

$$
\mathfrak{S}_{\chi}(a)=\mathfrak{S}_{\phi}(a), \quad \mathfrak{I}_{\chi}(a)=\mathfrak{T}_{\phi}(a)
$$

[The corresponding property of $\mathfrak{M}$ is included in Theorem 83.]
114. An increasing convex function of a convex function is convex.
115. If every chord of a continuous curve contains a point which lies above the curve, then every point of every chord, except the end-points, lies above the curve.
116. If $\phi(x)$ is convex and continuous, $a<b<c$, and $\phi(a)=\phi(b)=\phi(c)$, then $\phi(x)$ is constant in ( $a, c$ ).
117. If all the numbers are positive, then

$$
x \log \frac{x}{\bar{a}}+y \log \frac{y}{b}>(x+y) \log \frac{x+y}{a+b}
$$

unless $x / a=y / b$.
$\left[(x \log x)^{\prime \prime}>0\right.$.]
118. If $f(x)$ is positive and twice differentiable, then a necessary and sufficient condition that $\log f(x)$ should be convex is that $f f^{\prime \prime}-f^{\prime 2} \geqq 0$.
119. If $\phi(x)$ is continuous for $x>0$, and one of the functions $x \phi(x)$ and $\phi(1 / x)$ is convex, then so is the other.
120. If $\phi(x)$ is positive, twice differentiable and convex, then so are

$$
x^{\frac{1}{2}(s+1)} \phi\left(x^{-s}\right)(s \geqq 1), \quad e^{\frac{1}{2} x} \phi\left(e^{-x}\right)
$$

(the first for positive $x$ ).
121. If $\psi$ and $\chi$ are continuous and strictly monotonic, and $\chi$ increasing, then a necessary and sufficient condition that

$$
\psi^{-1}\left\{\psi^{a_{1}}\left(a_{1}\right) \ldots \psi^{q_{n}}\left(a_{n}\right)\right\} \leqq \chi^{-1}\left\{\chi^{q_{1}}\left(a_{1}\right) \ldots \chi^{a_{n}}\left(a_{n}\right)\right\}
$$

for all $a$ and $q$, is that

$$
\phi(y)=\log \left\{x \psi^{-1}\left(e^{y}\right)\right\}
$$

should be convex.
[Compare Theorem 92.]
122. Suppose that
(i)

$$
\phi\left(x_{1}\right)\left(x_{3}-x_{2}\right)+\phi\left(x_{2}\right)\left(x_{1}-x_{3}\right)+\phi\left(x_{3}\right)\left(x_{2}-x_{1}\right) \geqq 0,
$$

or (what is the same thing)
(ii)

$$
\left|\begin{array}{lll}
1 & x_{1} & \phi\left(x_{1}\right) \\
1 & x_{2} & \phi\left(x_{2}\right) \\
1 & x_{3} & \phi\left(x_{3}\right)
\end{array}\right| \geqq 0,
$$

for all $x_{1}, x_{2}, x_{3}$ of an open interval $I$ for which $x_{1}<x_{2}<x_{3}$. Then $\phi(x)$ is continuous in $I$, and has finite left-handed and right-handed derivatives at every point of $I$.

If $\phi(x)$ is twice differentiable, then (i) and (ii) are equivalent to the differential inequality

$$
\phi^{\prime \prime}(x) \geqq 0 .
$$

[(i) and (ii) are alternative forms of (3.5.2), and $\phi(x)$ is convex, so that the theorem is a restatement of parts of Theorems 111 and 94.]
123. Suppose that
(i)

$$
\phi\left(x_{1}\right) \sin \left(x_{3}-x_{2}\right)+\phi\left(x_{2}\right) \sin \left(x_{1}-x_{3}\right)+\phi\left(x_{3}\right) \sin \left(x_{2}-x_{1}\right) \geqq 0,
$$

or (what is the same thing)
(ii)

$$
\left|\begin{array}{l}
\cos x_{1} \sin x_{1} \phi\left(x_{1}\right) \\
\cos x_{2} \sin x_{2} \phi\left(x_{2}\right) \\
\cos x_{3} \sin x_{3} \phi\left(x_{3}\right)
\end{array}\right| \geqq 0,
$$

for all $x_{1}, x_{2}, x_{3}$ of an open interval $I$ for which $x_{1}<x_{2}<x_{3}<x_{1}+\pi$. Then $\phi(x)$ is continuous in $I$ and has finite left-handed and right-handed derivatives at every point of $I$.

If $\phi(x)$ is twice differentiable, then (i) and (ii) are equivalent to the differential inequality

$$
\phi^{\prime \prime}(x)+\phi(x) \geqq 0 .
$$

[The result is important in the study of convex curves and of the behaviour of analytic functions in angular domains. See Pólya (1, 320; 4, 573-576).]
124. A necessary and sufficient condition that a continuous function $\phi(x)$ should be convex in an interval $I$ is that, if $\alpha$ is any real number and $i$ any closed interval included in $I$, then $\phi(x)+\alpha x$ should attain its maximum in $i$ at one of the ends of $i$. If also $x$ and $\phi(x)$ are positive, then a necessary and sufficient condition that $\log \phi(x)$ should be a convex function of $\log x$ is that $x^{a} \phi(x)$ should have the same maximal property.
[For applications of this theorem, which results immediately from the definitions, see Saks (1).]
125. A necessary and sufficient condition that a continuous function $\phi(x)$ should be convex in $(a, b)$ is that
for

$$
\begin{gather*}
\phi(x) \leqq \frac{1}{2 h} \int_{x-h}^{x+h} \phi(t) d t  \tag{i}\\
a \leqq x-h<x<x+h \leqq b .
\end{gather*}
$$

[This is a corollary of Theorem 124. If $\phi(x)$ satisfies (i), so does $\phi(x)+\alpha x$; and it is plain that any continuous function which satisfies (i) must possess the property of Theorem 124. ${ }^{\text {a }}$

Theorem 125 may also be proved independently; and there are various generalisations. In particular we need only suppose (i) true for every $x$ and arbitrarily small $h=h(x)$.]

[^28]126. If $\phi(x)$ is convex and continuous for all $x$, and not constant, then $\phi(x)$ tends to infinity, for one or other approach of $x$ to infinity, in such a manner as to be ultimately greater than a constant multiple of $|x|$.
127. If $\phi^{\prime \prime}>0$ for $x>0$, and $\phi(0) \leqq 0$, then $\phi / x$ increases for $x>0$.
[This follows at once from the equations
$$
\left.x^{2} \frac{d}{d x}\left(\frac{\phi}{x}\right)=x \phi^{\prime}-\phi, \quad \frac{d}{d x}\left(x \phi^{\prime}-\phi\right)=x \phi^{\prime \prime} .\right]
$$
128. If $\phi^{\prime \prime}>0$ for $x \geqq 0$ and
$$
\lim _{x \rightarrow \infty}\left(x \phi^{\prime}-\phi\right) \leqq 0,
$$
then $\phi / x$ decreases (strictly) for $x>0$.
[The limit certainly exists, since $x \phi^{\prime}-\phi$ increases. The result follows from the equations used in proving Theorem 127. The cases considered in Theorems 127 and 128 are the extreme cases possible when $\phi^{\prime \prime}>0$; if neither condition is satisfied, $\phi / x$ has a minimum for some positive $x$.]
129. If $\phi^{\prime \prime}>0$ for all $x, \phi(0)=0$, and $\phi / x$ is interpreted as $\phi^{\prime}(0)$ for $x=0$, then $\phi / x$ increases for all $x$.
130. If the set $a_{1}, a_{2}, \ldots, a_{2 n+1}$ is convex in the sense of $\S 3.5$, i.e. if
then
\[

$$
\begin{gathered}
\Delta^{2} a_{\nu}=a_{\nu}-2 a_{\nu+1}+a_{\nu+2} \geqq 0 \quad(\nu=1,2, \ldots, 2 n-1), \\
\frac{a_{1}+a_{3}+\ldots+a_{2 n+1}}{n+1} \geqq \frac{a_{2}+a_{4}+\ldots+a_{2 n}}{n},
\end{gathered}
$$
\]

with inequality except when the numbers are in arithmetical progression.
[Nanson (1). Add up the inequalities

$$
r(n-r+1) \Delta^{2} a_{2 r-1} \geqq 0, \quad r(n-r) \Delta^{2} a_{2 r} \geqq 0 .
$$

Theorem 130 may also be proved as an example of Theorem 108: the set formed by the numbers $2,4, \ldots, 2 n$, each taken $n+1$ times, is majorised by that formed by the numbers $1,3, \ldots, 2 n+1$, each taken $n$ times.]
131. If $0<x<1$, then

$$
\frac{1-x^{n+1}}{n+1}>\frac{1-x^{n}}{n} \sqrt{ } x .
$$

[Put $x=y^{2}$ and $a_{\nu}=y^{\nu}$ in Theorem 130.]
132. $C$ is the centre of a circle and $A_{0} A_{1} \ldots A_{n} A_{0}$ an inscribed polygon, whose vertices lie on the circle in the order indicated. $C, A_{0}, A_{n}$ are fixed and $A_{1}, A_{2}, \ldots, A_{n-1}$ vary. Then the area and perimeter of the polygon are greatest when $\quad A_{0} A_{1}=A_{1} A_{2}=\ldots=A_{n-1} A_{n}$.
[Let $A_{\nu-1} C A_{\nu}=\alpha_{\nu}$. Since $(\sin x)^{\prime \prime}<0$ for $0<x<\pi$, we have, by Theorem 95,

$$
\frac{1}{n} \Sigma \sin \alpha_{\nu}<\sin \left(\frac{\Sigma \alpha_{\nu}}{n}\right)
$$

unless all the $\alpha_{\nu}$ are equal, and a similar inequality in which $\alpha_{\nu}$ is replaced by $\frac{1}{2} \alpha_{\nu}$. These inequalities give the two parts of the theorem. When $A_{n}$ coincides with $A_{0}$, they reduce to familiar maximal properties of regular polygons.]
133. Suppose that $f$ and $g$ are continuous increasing inverse functions which vanish at the origin, that $F=x f, G=x g$, and that $g$ satisfies the inequality

$$
g(x y) \leqq g(x) g(y)
$$

Suppose further that $\Sigma a b \leqq A B$ for all positive $a$ and $b$ such that $\Sigma G(b) \leqq G(B)$. Then

$$
\Sigma F(a) \leqq \frac{1}{F^{\prime}(1 / A)} .
$$

[Cooper (3). The result is included in Theorem 15 when $f$ is a power of $x$.]
134. If $\phi(x)$ is convex and continuous for $x \geqq 0, \nu=1,2,3, \ldots$, and the $a_{\nu}$ are non-negative and decreasing, then

$$
\phi(0)+\Sigma\left\{\phi\left(n a_{n}\right)-\phi\left[(n-1) a_{n}\right]\right\} \leqq \phi\left(\Sigma a_{n}\right) .
$$

If also $\phi^{\prime}(x)$ is a strictly increasing function, there is equality only when the $a_{\nu}$ are equal up to a certain point and then zero.
[Hardy, Littlewood, and Pólya (2). Write

$$
s_{0}=0, \quad s_{\nu}=a_{1}+a_{2}+\ldots+a_{\nu} \quad(\nu \geqq 1)
$$

and

$$
\begin{gathered}
s_{\nu}+(\nu-1) a_{\nu}=s_{\nu-1}+\nu a_{\nu}=2 x \\
s_{\nu}-(\nu-1) a_{\nu}=2 h, \quad s_{\nu-1}-\nu a_{\nu}=2 h^{\prime} .
\end{gathered}
$$

It is easily verified that $\left|h^{\prime}\right| \leqq h$, with equality only if $a_{\nu}=0$ or

$$
a_{1}=a_{2}=\ldots=a_{\nu} .
$$

It follows from Theorem 110 that

$$
\phi\left(\nu a_{\nu}\right)-\phi\left\{(\nu-1) a_{\nu}\right\} \leqq \phi\left(s_{\nu}\right)-\phi\left(s_{\nu-1}\right),
$$

and the result follows by summation.]
135. If $q>1$ and $a_{\nu}$ decreases, then

$$
\Sigma\left\{\nu^{q}-(\nu-1)^{q}\right\} a_{\nu}^{q} \leqq\left(\Sigma a_{\nu}\right)^{q}
$$

[Example of Theorem 134.]
136. Suppose that $a$ is a function of $\nu_{1}, \nu_{2}, \ldots, \nu_{m}$; that $i_{1}, i_{2}, \ldots, i_{m}$ is a permutation of the numbers $1,2, \ldots, m$; and that the $\psi$ and $\chi$ are continuous and strictly monotonic and the $\chi$ increasing. Then, in order that

$$
\mathfrak{M}_{\psi_{m}}^{\nu_{m}} \ldots \mathfrak{M}_{\psi_{1}}^{\nu_{1}}(a) \leqq \mathfrak{M}_{\lambda_{i_{m}}}^{\nu_{i_{m}}} \ldots \mathfrak{M}_{\lambda_{i_{1}}}^{\nu_{i_{1}}}(a)
$$

for all $a$ and $q$, it is necessary that
(1) $\chi_{\mu}$ is convex with respect to $\psi_{\mu}$, for $\mu=1,2, \ldots, m$;
(2) $\chi_{\lambda}$ is convex with respect to $\psi_{\mu}$, if $\lambda>\mu$ and $\lambda$ and $\mu$ correspond to an inversion in the permutation $i_{1}, i_{2}, \ldots, i_{m}$ (i.e. if the order is $\ldots \mu, \ldots, \lambda, \ldots$ in the series $1,2, \ldots, m$ but $\ldots \lambda, \ldots, \mu, \ldots$ in $\left.i_{1}, i_{2}, \ldots, i_{m}\right)$.
[Jessen (3).]
137. In order that

$$
\mathfrak{M}_{r_{m}}^{\nu_{m}} \ldots \mathfrak{M}_{r_{1}}^{\nu_{1}}(a) \leqq \mathfrak{M}_{s_{i_{m}}}^{\nu_{i_{m}}} \ldots \mathfrak{M}_{s_{i_{1}}}^{\nu_{i_{1}}}(a)
$$

for all $a$ and $q$, it is necessary and sufficient that (1) $r_{\mu} \leqq s_{\mu}$ and (2) $r_{\mu} \leqq s_{\lambda}$ (the range of $\mu$ and $\lambda$ being defined as in Theorem 136).
[Jessen (3). The most important cases are:
(i) $m=1, r<s$.
(ii) $m=2,\left(i_{1}, i_{2}\right)=(2,1), s_{2}=r_{2} \geqq s_{1}=r_{1}$.

The two cases correspond to Theorems 16 and 26. The kernel of the theorem is contained in the statement that, whenever the two sides of the inequality are comparable, the inequality may be proved by repeated application of the special cases corresponding to (i) and (ii).]
138. A continuous curve $y=\phi(x)$ defined in an open interval, say $0<x<1$, has the following property: through every point of the curve there is either ( $a$ ) a line which lies under the curve, or (b) a line which lies over the curve. Then one and the same of $(a)$ and $(b)$ is true at all points of the curve, and the curve is convex or concave.
[It is easy to show that if $S_{a}$ and $S_{b}$ are the sets of values of $x$ for which (a) or (b) is true, then $S_{a}$ and $S_{b}$ are closed (in the open interval). But a continuum cannot be the sum of two non-nul closed and exclusive sets.]
139. Suppose that $\phi(x)$ is convex and bounded below in ( $H, K$ ), and that $m(x)$ is the lower bound of $\phi(x)$ at $x$ (the limit of the lower bound of $\phi(x)$ in an interval including $x$ ). Then $m(x)$ is a continuous convex function; and either (i) $\phi(x)$ is identical with $m(x)$, or (ii) the graph of $\phi(x)$ is dense in the region $H \leqq x \leqq K, y \geqq m(x)$.

If $\phi(x)$ is convex and not bounded below, its graph is dense in the strip $H \leqq x \leqq K$.
[Bernstein and Doetsch (1).]

## CHAPTER IV

## VARIOUS APPLICATIONS OF THE CALCULUS

4.1. Introduction. Particular inequalities arising in ordinary analysis are often proved more easily by some special device than by an appeal to any general theory. We therefore interrupt our systematic treatment of the subject at this point, and devote a short chapter to the illustration of the simplest and most useful of these devices. The subject-matter is arranged according to the methods and instruments used rather than the character of the results.

## DIFFERENTIAL CALCULUS: FUNCTIONS OF ONE VARIABLE

4.2. Applications of the mean value theorem. Our first examples depend upon a straightforward use of the mean value theorem

$$
\begin{equation*}
f(x+h)-f(x)=h f^{\prime}(x+\theta h) \quad(0<\theta<1) \tag{4.2.1}
\end{equation*}
$$

or its generalisations with higher derivatives. It is a corollary of (4.2.1) that a function with a positive differential coefficient increases with $x$.
(1) We have

$$
\log (x+1)-\log x=\frac{1}{\xi},
$$

where $x<\xi<x+1$, when $x>0$. It follows that

$$
\begin{aligned}
& \frac{d}{d x}[x\{\log (x+1)-\log x\}]=\log (x+1)-\log x-\frac{1}{x+1}>0, \\
& \frac{d}{d x}[(x+1)\{\log (x+1)-\log x\}]=\log (x+1)-\log x-\frac{1}{x}<0 .
\end{aligned}
$$

Hence $\left(1+\frac{1}{x}\right)^{x}$ increases with $x$, while $\left(1+\frac{1}{x}\right)^{x+1}$ decreases.

Since the latter function is $\left(1-\frac{1}{y}\right)^{-y}$, where $y=x+1>1$, we obtain
140. $\left(1+\frac{1}{x}\right)^{x}$ increases for $x>0 ;\left(1-\frac{1}{x}\right)^{-x}$ decreases for $x>1$.

This is substantially the same as Theorem 35.
(2) If $x>1, r>1$, we have

$$
x^{r}=1+r(x-1)+\frac{1}{2} r(r-1) \xi^{r-2}(x-1)^{2},
$$

where $1<\xi<x$, and so
141. $x^{r}>1+r(x-1)+\frac{1}{2} r(r-1)\left(\frac{x-1}{x}\right)^{2} \quad(x>1, r>1)$.

This inequality was found, in a less precise form, in §2.15.
(3) If $x \neq 0$ we have

$$
\begin{equation*}
e^{x}=1+x+\frac{1}{2} x^{2} e^{\theta x}, \tag{4.2.2}
\end{equation*}
$$

where $0<\theta<1$, and so
142.

$$
e^{x}>1+x \quad(x \neq 0) .
$$

We can deduce another proof of Theorem 9. If

$$
\Sigma q=1, \quad \Sigma q a=\mathfrak{A}
$$

and the $a$ are not all equal, we can write $a=(1+x) \mathfrak{U}$, where $\Sigma q x=0$ and the $x$ are not all zero. Then $1+x \leqq e^{x}$, with inequality for at least one $x$, and

$$
\Pi a^{q}=\mathfrak{A} \Pi(1+x)^{q}<\mathfrak{A} e^{\Sigma q x}=\mathfrak{A}=\Sigma q a .
$$

The argument is a special case of that used at the end of §3.19.
(4) The function $f(x)=e^{x}-1-x-\frac{1}{2} x^{2}$,
and its first two derivatives, vanish for $x=0$. There is no other zero of $f(x)$, since this would (by repeated application of Rolle's Theorem) involve the existence of a zero of $f^{\prime \prime \prime}(x)=e^{x}$. Hence

$$
e^{x}>1+x+\frac{1}{2} x^{2} \quad(x>0), \quad e^{x}<1+x+\frac{1}{2} x^{2} \quad(x<0) .
$$

The same argument may be applied to any number of terms of the Taylor series of various functions. When the function is $e^{x}$, we obtain
143. If $n$ is odd then

$$
\begin{equation*}
e^{x}>1+x+\frac{x^{2}}{2!}+\ldots+\frac{x^{n}}{n!} \quad(x \neq 0) \tag{4.2.3}
\end{equation*}
$$

If $n$ is even then (4.2.3) is true for $x>0$ and the reversed inequality for $x<0$.
4.3. Further applications of elementary differential calculus. In this section we give some applications of a less immediate character.
(1) Repeated application of Rolle's Theorem leads easily to the following lemma ${ }^{\text {a }}$ : if

$$
f(x, y)=c_{0} x^{m}+c_{1} x^{m-1} y+\ldots+c_{m} y^{m}=0
$$

has all its roots $x / y$ real, then the same is true of all non-identical ${ }^{b}$ equations obtained from it by partial differentiations with respect to $x$ and $y$. Further, if $E$ is any one of these equations, and has a multiple root $\alpha$, then $\alpha$ is also a root, of multiplicity one higher, of the equation from which $E$ is derived by differentiation.

We use this lemma to prove a theorem proved already in a less complete form in § 2.22.
144. ${ }^{\mathrm{c}}$ If $a_{1}, a_{2}, \ldots, a_{n}$ are $n$ real, positive or negative, numbers, $p_{0}=1$, and $p_{\mu}$ is the arithmetic mean of the products of $\mu$ different $a$, then

$$
p_{\mu-1} p_{\mu+1}<p_{\mu}^{2} \quad(\mu=1,2, \ldots, n-1)
$$

unless all the a are equal.
We suppose that no $a$ is zero, since the specification of the cases of equality becomes more troublesome when zero $a$ 's are admitted.

Let $\quad f(x, y)=\left(x+a_{1} y\right)\left(x+a_{2} y\right) \ldots\left(x+a_{n} y\right)$

$$
=p_{0} x^{n}+\binom{n}{1} p_{1} x^{n-1} y+\binom{n}{2} p_{2} x^{n-2} y^{2}+\ldots+p_{n} y^{n} .
$$

Since no $a$ vanishes, $p_{n} \neq 0$ and $x / y=0$ is not a root of $f=0$. Hence $x / y=0$ cannot be a multiple root of any of the derived

[^29]equations; and therefore no two consecutive $p$, such as $p_{\mu}$ and $p_{\mu+1}$, can vanish. Hence the equation
$$
p_{\mu-1} x^{2}+2 p_{\mu} x y+p_{\mu+1} y^{2}=0
$$
which is obtained from $f(x, y)=0$ by a series of differentiations, is not identical; and therefore its roots are real, so that
$$
p_{\mu-1} p_{\mu+1} \leqq p_{\mu}^{2}
$$

Finally, the roots of the derived equation can be equal only if all the roots of the original equation are equal.

It will be observed that the $a$ need not necessarily be positive, as they were in § 2.22. ${ }^{\text {a }}$
(2) Suppose that $\phi(x)=\log \left(\Sigma p a^{x}\right)$, and (what is no real limitation) that the $a$ are all positive and unequal. Then

$$
\phi^{\prime}=\frac{\Sigma p a^{x} \log a}{\Sigma p a^{x}}, \quad \phi^{\prime \prime}=\frac{\Sigma p a^{x} \Sigma p a^{x}(\log a)^{2}-\left(\Sigma p a^{x} \log a\right)^{2}}{\left(\Sigma p a^{x}\right)^{2}}>0
$$

(Theorem 7). An easy calculation shows that, if $a_{r}$ is the greatest $a$, then

$$
\phi(0)=\log \Sigma p, \quad \lim _{x \rightarrow \infty}\left(x \phi^{\prime}-\phi\right)=-\log p_{r} .
$$

It follows from Theorems 127 and 129 that $\phi / x$ increases for $x>0$ if $\Sigma p \leqq 1$, and for all $x$ if $\Sigma p=1$. In the last case

$$
\frac{\phi(x)}{x}=\log \mathfrak{M}_{x}(a), \quad \lim _{x \rightarrow 0} \frac{\phi(x)}{x}=\phi^{\prime}(0)=\log (\mathscr{S}(a) .
$$

We thus obtain further proofs of Theorems 9 and 16.
If, on the other hand, $p_{\nu} \geqq 1$ for every $\nu, \phi / x$ decreases, by Theorem 128. In particular $\Im_{x}(a)=\left(\Sigma a^{x}\right)^{1 / x}$ decreases (Theorem 19). The general case gives part of Theorem 23.
(3) The following examples have applications in ballistics.
145.

$$
\log \sec x<\frac{1}{2} \sin x \tan x \quad\left(0<x<\frac{1}{2} \pi\right)
$$

146. The function

$$
\rho(x)=\frac{8 \log \sec x}{(g(x))^{2}}
$$

where

$$
g(x)=\int_{0}^{x}(1+\sec t) d t=x+\log (\sec x+\tan x),
$$

decreases steadily from 1 to 0 as $x$ increases from 0 to $\frac{1}{2} \pi$.
[Use Theorem 145 to show that $\frac{d}{d x}\left(g^{3} \rho^{\prime} \cot x\right)<0$ and so $\rho^{\prime}<0$.]
${ }^{\text {a }}$ Except in Theorem 55. For positive $a$, see Theorems 51, 54, 77, and §3.5.
147. The function

$$
\sigma(x)=\frac{\int_{0}^{x}(1+\sec t) \log \sec t d t}{. \log \sec x \int_{0}^{x}(1+\sec t) d t}
$$

increases steadily from $\frac{1}{3}$ to $\frac{1}{2}$ as $x$ increases from 0 to $\frac{1}{2} \pi$.
There is a general theorem which will be found useful in the proof of Theorem 147.
148. If $f, g$, and $f^{\prime} / g^{\prime}$ are positive increasing functions, then $f / g$ either increases for all $x$ in question, or decreases for all such $x$, or decreases to a minimum and then increases. In particular, if $f(0)=g(0)=0$, then $f / g$ increases for $x>0$.

To prove this, observe that

$$
\frac{d}{d x}\left(\frac{f}{g}\right)=\left(\frac{f^{\prime}}{g^{\prime}}-\frac{f}{g}\right) \frac{g^{\prime}}{g}
$$

and consider the possible intersections of the curves $y=f / g, y=f^{\prime} / g^{\prime}$. At one of these intersections the first curve has a horizontal and the second a rising tangent, and therefore there can be at most one intersection.
If we take $g$ as the independent variable, write $f(x)=\phi(g)$, and suppose, as in the last clause of the theorem, that

$$
f(0)=g(0)=0,
$$

or $\phi(0)=0$, then the theorem takes the form: if $\phi(0)=0$ and $\phi^{\prime}(g)$ increases for $g>0$, then $\phi / g$ increases for $g>0$. This is a slight generalisation of part of Theorem 127 . Theorem 148 should also be compared with Theorems 128 and 129.
4.4. Maxima and minima of functions of one variable. A very common method for the proof of inequalities is that of finding the absolute maximum or minimum of a function $\phi(x)$ by a discussion of the sign of $\phi^{\prime}(x)$.
(1) Since

$$
\frac{d}{d x}\left\{(1-x) e^{x}\right\}=-x e^{x}
$$

the function $(1-x) e^{x}$ has just one maximum, for $x=0$. Hence
149.

$$
e^{x}<\frac{1}{1-x} \quad(x<1, x \neq 0) .
$$

This is also a consequence of Theorem 142.
(2) Since

$$
\frac{d}{d x}(\log x-x+1)=\frac{1}{x}-1
$$

the function $\log x-x+1$ has one maximum, for $x=1$. Hence
150.

$$
\log x<x-1 \quad(x>0, x \neq 1)
$$

More generally $\log x<n\left(x^{1 / n}-1\right)$, for any positive $n$, since we may write $x^{1 / n}$ for $x$ in Theorem 150. This result is also a corollary of Theorem 36.
(3) Let $\phi(x)=1+x y-\left(1+x^{k}\right)^{1 / k}\left(1+y^{k^{\prime}}\right)^{1 / k^{\prime}}$,
where $k>1, x>0, y>0$. It is easily verified that $\phi(x)$ has a unique maximum 0 for $x^{k}=y^{k^{\prime}}$.

This gives another proof of $\mathrm{H}_{0}$ (Theorem 38), and so of H (Theorem 11).
(4) If $x$ and $y$ are positive, and $k>1$, then the function

$$
\phi(x)=x y-\frac{x^{k}}{k}-\frac{y^{k^{\prime}}}{k^{\prime}}
$$

has the derivative $y-x^{k / k^{\prime}}$, and attains its maximum 0 for $x^{k}=y^{k^{\prime}}$. We deduce Theorem 61 (and so Theorems 37 and 9).
(5) The function

$$
\phi(x)=x y-x \log x-e^{y-1},
$$

where $x$ is positive, attains its maximum 0 for $x=e^{y-1}$. We deduce Theorem 63.
4.5. Use of Taylor's series. If $f(x)=\Sigma a_{n} x^{n}$ and $g(x)=\Sigma b_{n} x^{n}$ are two series with positive coefficients, and $a_{n} \leqq b_{n}$ for every $n$, we say that $f(x)$ is majorised by $g(x)$, and write $f<g$. It is plain that $f \prec g$ and $f_{1} \prec g_{1}$ imply $f f_{1} \prec g g_{1}$, and so on.

To illustrate the use of this idea in the proof of inequalities, we prove
151. If $s_{n}=a_{1}+a_{2}+\ldots+a_{n}$, where $n>1$ and the $a$ are positive, then

$$
\left(1+a_{1}\right)\left(1+a_{2}\right) \ldots\left(1+a_{n}\right)<1+\frac{s_{n}}{1!}+\frac{s_{n}{ }^{2}}{2!}+\ldots+\frac{s_{n}{ }^{n}}{n!} .
$$

In fact $1+a_{\nu} x<e^{a_{\nu} x}$, so that

$$
\Pi(1+a x)<e^{s_{n} x} .
$$

The result follows by adding up the coefficients of $1, x, x^{2}, \ldots, x^{n}$, and observing that there is strict inequality between the coefficients of $x^{2}$. It may also be proved by writing the left-hand side in the form

$$
1+n p_{1}+\frac{n(n-1)}{1.2} p_{2}+\ldots+p_{n}
$$

(so that $n p_{1}=s_{n}$ ) and using Theorem 52.

## DIFFERENTIAL CALCULUS: FUNCTIONS OF SEVERAL VARIABLES

4.6. Applications of the theory of maxima and minima of functions of several variables. The most 'universal' weapon for the discovery and proof of inequalities is the general theory of maxima and minima of functions of any number of variables. Suppose that we wish to prove that two functions $\phi$ and $\psi$ of the continuous variables $x_{1}, x_{2}, \ldots, x_{n}$ are comparable; let us say, to fix ourideas, that $\phi-\psi \geqq 0$. This will be so, if $\phi-\psi$ has a minimum, if and only if this minimum is non-negative; and this is a question which can always be attacked (at any rate when the functions are differentiable) by the standard arguments of the theory of maxima and minima.

The method is attractive theoretically, and always opens a first line of attack on the problem; but it is apt to lead to serious complications in detail (usually connected with the 'boundary values' of the variables), and it will be found that, however suggestive it may be, it rarely leads to the simplest solution. We illustrate these remarks by considering its application to the fundamental inequalities G and H .
(1) To prove G, consider
where

$$
\begin{aligned}
& f\left(x_{1}, x_{2}, \ldots, x_{n-1}\right)=x_{1}{ }_{q_{1}} x_{2}^{q_{2}} \ldots x_{n}^{q_{n}}, \\
& x_{n}=\frac{1}{q_{n}}\left(\mathfrak{A}-q_{1} x_{1}-\ldots-q_{n-1} x_{n-1}\right),
\end{aligned}
$$

in the closed and bounded domain $x_{1} \geqq 0, \ldots, x_{n} \geqq 0$. It is continuous, and therefore attains a maximum, which is not on the boundary (where $f$ vanishes). At the maximum

$$
0=\frac{1}{f} \frac{\partial f}{\partial x_{v}}=\frac{q_{\nu}}{x_{\nu}}-\frac{q_{n}}{x_{n}} \frac{q_{\nu}}{q_{n}} \quad(\nu=1,2, \ldots, n-1),
$$

and so the $x$ are all equal to $\mathfrak{A}$. In this case no serious complication is introduced by the boundary values ${ }^{\text {a }}$.
(2) We may use H (for two sets of variables) as an illustration of the 'method of Lagrange'. Consider

$$
\begin{gathered}
f\left(x_{1}, x_{2}, \ldots, x_{n}\right)=b_{1} x_{1}+b_{2} x_{2}+\ldots+b_{n} x_{n}, \\
\text { a Compare } \S 2.6(\mathrm{i}) .
\end{gathered}
$$

where $b_{\nu}>0$, subject to the condition that

$$
\phi\left(x_{1}, x_{2}, \ldots, x_{n}\right)=x_{1}^{k}+x_{2}^{k}+\ldots+x_{n}^{k} \quad(k>1)
$$

is a positive constant $X$. The ( $n-1$ )-dimensional domain defined by $x \geqq 0, \phi=X$ is closed and bounded, and some $x$ vanishes at every point of its boundary.

If the maximum is attained at an internal point, then, at that point,

$$
\frac{\phi_{x_{\nu}}}{f_{x_{\nu}}}=\frac{k x_{\nu}^{k-1}}{b_{\nu}}=\lambda
$$

is independent of $\nu$; and an elementary calculation shows that

$$
f=X^{1 / k}\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}}=\left(\Sigma x^{k}\right)^{1 / k}\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}} .
$$

There remains the possibility that the maximum should be attained at a boundary point, where some $x$, say $x_{n}$, is zero. This possibility may be excluded by an inductive argument, since, if we assume that the inequality has been proved already for $n-1$ variables, and that $x_{n}=0$, we have

$$
\left.f=\sum_{1}^{n-1} b_{\nu} x_{\nu} \leqq\left(\sum_{1}^{n-1} x_{\nu}\right)^{k}\right)^{1 / k}\left(\sum_{1}^{n-1} b_{\nu} k^{\prime}\right)^{1 / k^{\prime}}<\left(\Sigma x_{\nu}{ }^{k}\right)^{1 / k}\left(\Sigma b_{\nu}^{k^{\prime}}\right)^{1 / k^{\prime}}
$$

The weakness of the method is that, if we are to argue by induction at all, it is better to prove the whole theorem inductively, and then we come back to one of the proofs of H given already.
(3) It is quite usual that the method should, as in case (2), prove troublesome when developed in detail; but even in such cases it is very useful as indicating a possible solution of the problem.

A great many of our theorems assert inequalities between two symmetric functions $f\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ and $g\left(x_{1}, x_{2}, \ldots, x_{n}\right)$, homogeneous of the same degree and positive for all positive $x$. This is true, for example, of Theorems 9,16 and 17 (for unit weights, the crucial case), Theorem 45 (in the case of comparability) and Theorems 51 and 52.

When we use Lagrange's method, we must consider the maximum of $f$ for constant $g$, say for $g=1$. Lagrange's equations are

$$
\frac{\partial f}{\partial x_{\nu}}=\lambda \frac{\partial g}{\partial x_{\nu}} \quad(\nu=1,2, \ldots, n)
$$

These equations have always a solution with $x_{1}=x_{2}=\ldots=x_{n}$, and $\lambda$ is the value of $f$ for this system of values of $x$. If $\lambda$ is a strict absolute maximum of $f$, then $f \leqq \lambda g$, and there is inequality except when all the $x$ are equal.

All this is in fact true in the cases mentioned, but there are other cases in which the solution does not give the maximum of $f$. This happens for example in Theorem 45, in the case of non-comparability.

## INTEGRAL CALCULUS

4.7. Comparison of series and integrals. There are many inequalities which are proved most easily by arguments based on the integral calculus; and often, by consideration of areas or volumes, in an 'intuitive' way. We give here a few of the most useful general theorems, in which the integrals considered are of the elementary Riemann or Riemann-Stieltjes type. In Ch. VI we shall consider inequalities between integrals systematically, and there we shall use the general Lebesgue and LebesgueStieltjes integrals.

The theorems which follow immediately are due in principle to Maclaurin and Cauchy ${ }^{\text {a }}$.
152. If $f(x)$ decreases for $x \geqq 0$, then

$$
f(1)+f(2)+\ldots+f(n) \leqq \int_{0}^{n} f(x) d x \leqq f(0)+f(1)+\ldots+f(n-1)
$$

There is inequality if $f(x)$ decreases strictly.
In fact

$$
f(\nu+1) \leqq \int_{\nu}^{\nu+1} f(x) d x \leqq f(\nu)
$$

(with inequality if $f(x)$ decreases strictly).
Further theorems of the same type, which we state without proof, are:
153. If $a_{0}<a_{1}<a_{2}<\ldots$, and $f(x)$ decreases for $x \geqq a_{0}$, then

$$
\begin{gathered}
\sum_{\nu=1}^{n}\left(a_{\nu}-a_{\nu-1}\right) f\left(a_{\nu}\right) \leqq \int_{a_{0}}^{a_{n}} f(x) d x \leqq \sum_{\nu=1}^{n}\left(a_{\nu}-a_{\nu-1}\right) f\left(a_{\nu-1}\right) . \\
\text { a } \operatorname{Maclaurin}(1,1,289) ; \operatorname{Cauchy}(2,222) .
\end{gathered}
$$

154. If $f(x) \geqq 0$, and $f(x)$ decreases in $(0, \xi)$ and increases in $(\xi, 1)$, where $0 \leqq \xi \leqq 1$, then

$$
\begin{aligned}
\frac{1}{n}\left\{f\left(\frac{1}{n}\right)+f\left(\frac{2}{n}\right)+\ldots+f\left(\frac{n-1}{n}\right)\right\} & \leqq \int_{0}^{1} f(x) d x \\
& \leqq \frac{1}{n}\left\{f(0)+f\left(\frac{1}{n}\right)+\ldots+f(1)\right\}
\end{aligned}
$$

155. If $f(x, y)$ is a decreasing function of $x$ for fixed $y$ and a decreasing function of $y$ for fixed $x$, then

$$
\sum_{\mu=1}^{m} \sum_{\nu=1}^{n} f(\mu, \nu) \leqq \int_{0}^{m} \int_{0}^{n} f(x, y) d x d y \leqq \sum_{\mu=0}^{n-1} \sum_{\nu=0}^{n-1} f(\mu, \nu)
$$

Applications of these theorems, particularly to the theory of the convergence of series, may be found in any textbook of analysis.
4.8. An inequality of Young. The simple but useful theorem which follows is due to W.H. Young ${ }^{\text {a }}$ and is of a different type.
156. Suppose that $\phi(0)=0$, and that $\phi(x)$ is continuous and strictly increasing for $x \geqq 0$; that $\psi(x)$ is the inverse function, so that $\psi(x)$ satisfies the same conditions; and that $a \geqq 0, b \geqq 0$. Then

$$
a b \leqq \int_{0}^{a} \phi(x) d x+\int_{0}^{b} \psi(x) d x .
$$

There is equality only if $b=\phi(a)$.
The theorem becomes intuitive if we draw the curve $y=\phi(x)$ or $x=\psi(y)$, and the lines $x=0, x=a, y=0, y=b$, and consider the various areas bounded by them. A formal proof is included in that of the more general theorems which follow.

A corollary of Theorem 156 is
157. If the conditions of Theorem 156 are satisfied, then

$$
a b \leqq a \phi(a)+b \psi(b)
$$

Theorem 157 is weaker than Theorem 156, but often aseffective in applications.

We pass to more general theorems which include Theorem 156.

[^30]158. Suppose that $\nu=1,2, \ldots, n$; that $a_{\nu} \geqq 0$; that $f_{v}(x)$ is continuous, non-negative, and strictly increasing; and that one of the $f_{\nu}(0)$ is zero. Then
$$
\Pi f_{\nu}\left(a_{\nu}\right) \leqq \Sigma \int_{0}^{a_{\nu}} \Pi_{\mu \neq \nu} f_{\mu}(x) . d f_{\nu}(x)
$$
and there is equality only if $a_{1}=a_{2}=\ldots=a_{n} \cdot{ }^{\text {a }}$
The inequality may be made intuitive by considering the curve $x_{\nu}=f_{\nu}(t)$ in $n$-dimensional space, and the volumes bounded by the coordinate planes and the cylinders which project the curves on these planes.

To obtain a formal proof, put

$$
F_{\nu}(x)=f_{\nu}(x) \quad\left(0 \leqq x \leqq a_{\nu}\right), \quad F_{\nu}(x)=f_{\nu}\left(a_{\nu}\right) \quad\left(x \geqq a_{\nu}\right),
$$

so that $F_{\nu}(x) \leqq f_{\nu}(x)$. If we suppose, as we may, that $a_{n}$ is the largest of the $a_{\nu}$, then, since $\Pi F_{\nu}(0)=0$, we have

$$
\begin{aligned}
\Pi f_{\nu}\left(a_{\nu}\right)=\Pi F_{\nu}\left(a_{n}\right) & =\int_{0}^{a_{n}} d\left\{\Pi F_{\nu}(x)\right\}=\sum_{\nu} \int_{0}^{a_{n}} \prod_{\mu \neq \nu} F_{\mu}(x) \cdot d F_{\nu}(x) \\
& =\sum_{\nu} \int_{0}^{a_{\nu}} \Pi_{\mu \neq \nu} F_{\mu}(x) \cdot d f_{\nu}(x) \leqq \sum_{\nu}^{a_{\nu}} \int_{0}^{a_{\mu} \neq \nu} f_{\mu}(x) \cdot d f_{\nu}(x)
\end{aligned}
$$

and there is inequality unless every $a_{\nu}$ is equal to $a_{n}$.
159.b Suppose that $g_{\nu}(x)$ is a system of $n$ continuous and strictly increasing functions each of which vanishes for $x=0$; that

$$
\begin{equation*}
\Pi g_{\nu}^{-1}(x)=x \tag{4.8.1}
\end{equation*}
$$

and that $a_{\nu} \geqq 0$. Then

$$
\Pi a_{\nu} \leqq \Sigma \int_{0}^{a_{\nu}} \frac{g_{\nu}(x)}{x} d x
$$

There is inequality unless $g_{1}\left(a_{1}\right)=g_{2}\left(a_{2}\right)=\ldots=g_{n}\left(a_{n}\right)$.
If we put

$$
g_{\nu}^{-1}(x)=\chi_{\nu}(x), \quad b_{\nu}=g_{\nu}\left(a_{\nu}\right), \quad a_{\nu}=g_{\nu}^{-1}\left(b_{\nu}\right)=\chi_{\nu}\left(b_{\nu}\right)
$$

and apply Theorem 158 to the system $f_{\nu}(x)=\chi_{\nu}(x)$, we obtain

$$
\Pi a_{\nu}=\Pi \chi_{\nu}\left(b_{\nu}\right) \leqq \Sigma \int_{0}^{b_{\nu}} \frac{x}{\chi_{\nu}(x)} d \chi_{\nu}(x)=\Sigma \int_{0}^{a_{\nu}} \frac{g_{\nu}(y)}{y} d y
$$

a Oppenheim (1). The proof is by T. G. Cowling.
${ }^{\mathrm{b}}$ Cooper (1).

A system of $n$ functions

$$
\phi_{\nu}(x)=\frac{g_{\nu}(x)}{x}
$$

connected by (4.8.1) is a generalisation of a pair of inverse functions. For suppose that $n=2, g_{1}(x)=x \phi(x), g_{2}(x)=x \psi(x)_{2}$ and write (4.8.1) in the two forms

$$
\frac{x}{g_{1}^{-1}(x)}=g_{2}^{-1}(x), \quad \frac{x}{g_{2}^{-1}(x)}=g_{1}^{-1}(x) .
$$

Then

$$
\phi(x)=\frac{g_{1}(x)}{x}=g_{2}^{-1}\left\{g_{1}(x)\right\}, \quad \psi(x)=\frac{g_{2}(x)}{x}=g_{1}^{-1}\left\{g_{2}(x)\right\},
$$

and $g_{2}^{-1} g_{1}, g_{1}^{-1} g_{2}$ are inverse; and the functions $\phi$ and $\psi$ of Theorem 156 can always be represented in this form. Hence Theorem 159 includes Theorem 156.

If in Theorem 159 we take $g_{\nu}(x)=x^{1 / q_{\nu}}$, where $\Sigma q_{\nu}=1$, then (4.8.1) is satisfied, and we obtain

$$
\Pi a_{\nu} \leqq \Sigma q_{\nu} a_{\nu}^{1 / q_{\nu}}
$$

which is Theorem 9. If in Theorem 156 we take $\phi(x)=x^{k-1}$, where $k>1$, we have $\psi(x)=x^{k^{\prime}-1}$, and we obtain Theorem 61. If we take

$$
\phi(x)=\log (x+1), \quad \psi(y)=e^{y}-1,
$$

and write $u, v$ for $a+1$ and $b+1$, we obtain Theorem 63, for $u \geqq 1, v \geqq 1$. ${ }^{\mathrm{a}}$

We pointed out above that Theorem 156 is intuitive from simple graphical considerations. If instead of reckoning areas we count up the number of lattice-points (points with integral coordinates) inside them, we obtain
160. If $\phi(x)$ increases strictly with $x, \phi(0)=0$, and $\psi(x)$ is the function inverse to $\phi(x)$, then

$$
m n \leqq \sum_{0}^{m}[\phi(\mu)]+\sum_{0}^{n}[\psi(\nu)]
$$

where $[y]$ is the integral part of $y$.
This theorem is less interesting in itself, but illustrates a type of argument often effective in the Theory of Numbers.
a Actually the result is true for $u>0$ and all $v$. See §4.4(5).

## CHAPTER V

## INFINITE SERIES

5.1. Introduction. Our theorems so far have related to finite sums, and we have now to consider their extensions to infinite series. The general conclusion will be that our theorems remain valid for infinite series in so far as they retain significance.

Two preliminary remarks are necessary.
(1) The first concerns the interpretation of our formulae. An inequality $X<Y$ ( or $X \leqq Y$ ), where $X$ and $Y$ are infinite series, is always to be interpreted as meaning 'if $Y$ is convergent, then $X$ is convergent, and $X<Y$ (or $X \leqq Y$ ). More generally, an inequality of the type

$$
\begin{equation*}
X<\Sigma A Y^{b} \ldots Z^{c} \tag{5.1.1}
\end{equation*}
$$

(or $X \leqq \Sigma A Y^{b} \ldots Z^{c}$ ), where $Y, \ldots, Z$ are any finite number of infinite series, $\Sigma$ is a finite sum, and $A, b, \ldots, c$ are positive, is to be interpreted as meaning 'if $Y, \ldots, Z$ are convergent then $X$ is convergent, and $X$ satisfies the inequality'. Neglect of this understanding would lead to confusion when it is ' $<$ ' which stands in the inequality. We could read ' $Y$ ' as ' $\infty$ ' in the case of divergence; then ' $X \leqq \infty$ ' would convey no information, but ' $X<\infty$ ' would imply the convergence of $X$, and this implication would usually be false.

Some inequalities will occur which are not of the form (5.1.1). These are usually secondary, and should be reduced to the form (5.1.1) if there is any doubt about their interpretation. Thus

$$
X^{a}<A Y^{b}
$$

should be interpreted as

$$
X<A^{1 / a} Y^{b / a},
$$

which is of the form (5.1.1); and $X>Y$ should be interpreted as $Y<X$.

There is one important inequality, viz.

$$
\begin{equation*}
\Sigma a b>\left(\Sigma x^{k}\right)^{1 / k}\left(\Sigma b^{k^{\prime}}\right)^{1 / k^{\prime}}, \tag{5.1.2}
\end{equation*}
$$

where $k<1, k \neq 0$, a which we have written deliberately in a form unlike (5.1.1). We might have written it as

$$
\begin{equation*}
\Sigma a^{k}<(\Sigma a b)^{k}\left(\Sigma b^{k^{\prime}}\right)^{-k / k^{\prime}} \tag{5.1.3}
\end{equation*}
$$

when $0<k<1$, or as

$$
\begin{equation*}
\Sigma b^{k^{\prime}}<(\Sigma a b)^{k^{\prime}}\left(\Sigma a^{k}\right)^{-k^{\prime} / k} \tag{5.1.4}
\end{equation*}
$$

when $k<0$. These are of the form (5.1.1), and are the forms which arise primarily in the proof of Theorem 13. We prefer the form (5.1.2) for formal reasons, and because it shows clearly the contrast between the two cases of the theorem, but we must use the other forms if we wish to show explicitly and exactly the implications about convergence.

There are a few cases where the inequality asserted is not one between two infinite series but involves the results of other limit operations. Thus, when we extend the inequality ( $5(a)<\operatorname{Max} a$ (Theorem 2) we obtain an inequality between an infinite product and the upper bound of an infinite set. Such an inequality ' $X<Y$ ' is to be interpreted in the same way, viz. as 'if $Y$ is finite then $X$ is finite and $X<Y^{\prime}$.
(2) The second remark concerns method, and should be read in conjunction with §1.7. Suppose, for example, that we wish to prove the inequality $\quad(\Sigma a b)^{2} \leqq \Sigma a^{2} \Sigma b^{2}$
for infinite series. We know the inequality for finite sums (Theorem 7) and, everything being positive, our conclusion follows by a passage to the limit.

We cannot extend Theorem 7, in its complete form, to infinite series in so simple a manner, since in the limiting process ' $<$ ' degenerates into ' $\leqq$ ', and we are unable to pick out the possible cases of equality. Here and elsewhere, we must avoid limiting processes; instead of deducing the infinite theorem from the finite one, we must verify that the proof given for the finite theorem remains valid, with that minimum of change required by the new context, in the infinite case. For example, either proof of Theorem 7 given in §2.4 may be extended to the infinite

[^31]case by the addition of a few obvious comments concerning convergence.

It will not be necessary to retrace the path followed in Ch. II systematically. The few new points which arise are neither difficult nor particularly interesting, and, in so far as they are important, recur in a more interesting form in Ch. VI. We shall therefore arrange the substance of this chapter informally, illustrating and commenting upon the new possibilities, and ending with an enumeration of some of the more important theorems of Ch . II which remain valid with the new interpretation.
5.2. The means $\mathfrak{M}_{r}$. We begin by some comments on a new point which arises in the definition of the means $\mathfrak{M}_{r}$. We have now an infinity of terms $a$ and weights $p$, and there are two cases to consider, according as $\Sigma p$ is convergent or divergent.
(i) If $\Sigma p$ is convergent we may suppose that the sum is 1 and write $q$ for $p$. In this case $\mathfrak{M}_{r}$ is defined, for $r>0$, by

$$
\begin{equation*}
\mathfrak{M}_{r}(a)=\left(\Sigma q a^{r}\right)^{1 / r} \tag{5.2.1}
\end{equation*}
$$

and may be regarded as a 'mean' in the sense of $\S 2.2$ or a 'weighted sum' in the sense of $\S 2.10$ (iv). It is finite or infinite according as $\Sigma q a^{r}$ is convergent or divergent.
(ii) If $\Sigma p$ is divergent, we can still define $\mathfrak{M}_{r}$ as a limit, e.g. by

$$
\begin{equation*}
\mathfrak{M}_{r}(a)=\lim _{n \rightarrow \infty}\left(\sum_{1}^{n} p_{\nu} a_{\nu}^{r} / \sum_{1}^{n} p_{\nu}\right)^{1 / r} \tag{5.2.2}
\end{equation*}
$$

or as the corresponding upper limit $\varlimsup$. The latter definition is not particularly interesting, though it would preserve most of our theorems. If we define $\mathfrak{M}_{r}$ by (5.2.2), we are met by a difficulty: the existence of $\mathfrak{M}_{r}$ for a given $r$ does not ensure its existence for any other $r$. In fact we can determine the $a$ so that $M_{r}$ shall exist for a given set $r_{1}, r_{2}, \ldots, r_{m}$ of values of $r$ and for no others. We shall therefore confine our attention to case (i).

For the general question of the existence of $\mathfrak{M}_{r}$, see, for example, Besicovitch (1). We may illustrate the point by showing briefly how to find $a$ so that either of the limits

$$
\lim \frac{1}{n}\left(a_{1}+a_{2}+\ldots+a_{n}\right), \quad \lim \frac{1}{n}\left(a_{1}^{2}+a_{2}^{2}+\ldots+a_{n}^{2}\right)
$$

may exist without the other: here $p=1$.

Take first two sequences

$$
\alpha_{1}, \alpha_{2}, \ldots, \alpha_{w}, \alpha_{1}, \alpha_{2}, \ldots ; \beta_{1}, \beta_{2}, \ldots, \beta_{w}, \beta_{1}, \beta_{2}, \ldots
$$

with period $w$. When $a=\alpha$, both limits exist and have the values

$$
A_{1}=\frac{\alpha_{1}+\alpha_{2}+\ldots+\alpha_{w}}{w}, \quad A_{2}=\frac{\alpha_{1}{ }^{2}+\alpha_{2}{ }^{2}+\ldots+\alpha_{w}{ }^{2}}{w} ;
$$

and when $a=\beta$ they have the corresponding values $B_{1}, B_{2}$.
Now take the $a$ as follows:

It is easy to see that, by supposing the sequence $N_{1}, N_{2} / N_{1}, N_{3} / N_{2}, \ldots$ to tend to infinity with great rapidity, we can make ( $a_{1}+\ldots+a_{n}$ )/n and $\left(a_{1}^{2}+\ldots+a_{n}^{2}\right) / n$ oscillate between $A_{1}, B_{1}$ and $A_{2}, B_{2}$ respectively. The conditions for convergence will then be $A_{1}=B_{1}$ and $A_{2}=B_{2}$ respectively, and we can obviously choose the $\alpha$ and $\beta$ so that either of these conditions shall be satisfied without the other.

We therefore restrict ourselves to case (i). We define $\mathfrak{M}_{r}$, for positive or negative $r$, by (5.2.1), with the convention that $\mathfrak{M} \boldsymbol{M}_{r}=0$ if $r$ is negative and an $a$ zero or $\Sigma q a^{r}$ divergent. We define $\mathfrak{G}$ (and $\mathfrak{M}_{0}$ ) by

$$
\begin{equation*}
\mathscr{S}(a)=\mathfrak{M}_{0}(a)=\Pi a^{q}=\exp (\Sigma q \log a), \tag{5.2.3}
\end{equation*}
$$

with the conventions that $\mathbb{G}=\infty$ if $\Pi a^{q}$ diverges to $\infty$ (i.e. if $\Sigma q \log a$ diverges to $+\infty$ ), and $\mathscr{5}=0$ if $\Pi a^{q}$ diverges to 0 (i.e. if $\Sigma q \log a$ diverges to $-\infty)$. It is to be observed that $\log a$ may have either sign, and that the definition of $\mathscr{F}$ fails if $\Sigma q \log a$ is oscillatory. In this case $\mathfrak{F}$ is meaningless.

It follows from Theorems 36 and 150 that

$$
\begin{equation*}
\log t<\frac{t^{r}-1}{r}<\frac{t^{s}-1}{s} \quad(0<r<s, t>0), \tag{5.2.4}
\end{equation*}
$$

unless $t=1$, when there is equality. We define $\log ^{+} t$ and $\log ^{-} t$ by

$$
\begin{array}{ll}
\log ^{+} t=\log t(t>1), & \log ^{+} t=0 \quad(0<t \leqq 1), \\
\log -t=\log t(0<t \leqq 1), & \log ^{-} t=0 \quad(t>1),
\end{array}
$$

so that

$$
\log ^{+} t \geqq 0, \log ^{-} t \leqq 0, \log t=\log ^{+} t+\log ^{-} t, \log ^{-} t=-\log ^{+} \frac{1}{t}
$$

It then follows from (5.2.4) that

$$
0 \leqq \Sigma q \log ^{+} a \leqq \frac{1}{r} \Sigma^{\prime} q\left(a^{r}-1\right) \leqq \frac{1}{s} \Sigma^{\prime} q\left(a^{s}-1\right)
$$

where $\Sigma^{\prime}$ denotes a summation extended over the $a$ which exceed l. Hence, if $\mathbb{M}_{s}(a)$ is finite for some positive $s$, then $\mathbb{M}_{r}(a)$ is finite for $0<r<s$, and $\Sigma q \log ^{+} a$ is convergent. We can prove similarly that, if $\mathfrak{M}_{-s}(a)$ is positive for some positive $s$, then $\mathbb{M}_{-r}(a)$ is positive for $-s<-r<0$, and $\Sigma q \log -a$ is convergent. In the first case $(5)$ is positive and finite or zero, in the second it is positive and finite or infinite. If $\Sigma q \log a$ oscillates then $\Sigma q \log ^{+} a$ and $\Sigma q \log ^{-} a$ are both divergent, and this is only possible when $\mathfrak{M}_{r}(a)=\infty$ for all positive $r$ and $\mathfrak{M}_{r}(a)=0$ for all negative $r$. It is only in this case that the definition of $\mathscr{G}(a)$ can fail.

There is one new point which, as we shall see in §5.9, affects the specification of the cases of equality in some of our theorems. This point arises from the possibility that, when $r \leqq 0, \mathfrak{M}_{r}(a)$ may be zero although no $a$ is zero. If $r>0$ then, as in Ch. II, $\mathfrak{M}_{r}(a)$ can be zero only if $(a)$ is nul, in which case $\mathfrak{M}_{r}(a)=0$ for all $r$. But when $r \leqq 0$ there is a difference. The $\mathfrak{M r}_{r}(a)$ of Ch. II was zero, for such an $r$, if and only if some $a$ was zero, and then $\mathfrak{M}_{r}(a)$ was zero for all $r \leqq 0$. It is now possible, when $r \leqq 0$, that $\mathfrak{M}_{s}(a)$ should be zero for $s<r$ and positive for $s \geqq r$, or zero for $s \leqq r$ and positive for $s>r$.

Thus in Theorem 1 there were two exceptional cases;

$$
\operatorname{Min} a<\operatorname{Ma}_{r}(a)<\operatorname{Max} a
$$

unless either all the $a$ are equal, or else $r<0$ and an $a$ is zero. All that we can say now is 'unless either all the $a$ are equal (in which case both inequalities reduce to equalities) or else $r<0$ and $\mathfrak{M}_{r}(\alpha)=0$ (in which case $\operatorname{Min} a=0$ and the first inequality reduces to an equality)'. Substantially the same point arises in connection with Theorems 2, 5, 10, 16, 24 and 25 (to quote only cases referred to in our summary in §5.9).
5.3. The generalisation of Theorems 3 and 9. We use the inequalities (5.2.4), and the equation

$$
\lim _{r \rightarrow 0} \frac{t^{r}-1}{r}=\log t
$$

Taking $t=a / \Sigma q a=a / \mathfrak{A}, r=1$ in (5.2.4), we have

$$
\begin{gathered}
\log a-\log \mathfrak{A} \leqq \frac{a}{\mathfrak{\mathfrak { V }}}-1 \\
\log \mathfrak{G}-\log \mathfrak{A}=\Sigma q(\log a-\log \mathfrak{A}) \leqq 1-1=0
\end{gathered}
$$

with equality only if every $a$ is $\mathfrak{A}$. This proves the analogue of Theorem 9.

Suppose now that $\mathbb{M}_{s}$ is finite for some positive $s$. Then $\mathfrak{F s}$ is positive and finite or zero: the proof below applies to either case ${ }^{\text {a }}$. Given $\epsilon>0$, we can choose $N$ so that

$$
\begin{gather*}
\sum_{n \leqq N} q \log a<\log (\mathscr{S}+\epsilon),  \tag{5.3.1}\\
\sum_{n>N} q \frac{a^{s}-1}{s}<\epsilon,
\end{gather*}
$$

and then $r_{0}$ so that $0<r_{0}<s$ and

$$
\begin{equation*}
\sum_{n \leqq N} q \frac{a^{r}-1}{r}<\sum_{n \leqq N} q \log a+\epsilon \tag{5.3.3}
\end{equation*}
$$

for $0<r<r_{0}$. We have then

$$
\begin{aligned}
& \log \mathfrak{S}(a)=\frac{1}{r} \log \left(\mathscr{S}\left(a^{r}\right) \leqq \frac{1}{r} \log \mathfrak{A}\left(a^{r}\right) \leqq \frac{\mathfrak{P}\left(a^{r}\right)-1}{r}\right. \\
& =\sum_{n \leqq N} q \frac{a^{r}-1}{r}+\sum_{n>N} q \frac{a^{r}-1}{r}<\sum_{n \leqq N} q \log a+\epsilon+\sum_{n>N} q \frac{a^{s}-1}{s} \\
& <\log (\mathscr{S}+\epsilon)+2 \epsilon \text {. } \\
& \log \mathfrak{M}_{r}(a)=\frac{1}{r} \log \mathfrak{U}\left(a^{r}\right) \rightarrow \log \mathscr{S}(a) \\
& \text { Hence }
\end{aligned}
$$

when $r \rightarrow+0$. We leave it to the reader to prove that, if $\mathbb{M}_{r}$ is positive for some negative $r$, then $\mathfrak{M}_{r} \rightarrow \mathfrak{F}$ when $r \rightarrow-0$.
5.4. Hölder's inequality and its extensions. The proofs of Hölder's inequality, and other theorems of the same type, given in Ch. II apply equally to infinite series. We may observe in passing that the series may be multiple series. Thus

$$
\left(\Sigma \Sigma a_{\mu \nu} b_{\mu \nu}\right)^{2} \leqq \Sigma \Sigma a_{\mu \nu}^{2} \Sigma \Sigma b_{\mu \nu}^{2}
$$

[^32]For example, suppose that $\Sigma u_{\mu}{ }^{2}$ and $\Sigma v_{\nu}{ }^{2}$ are convergent, and take

$$
a_{\mu \nu}=u_{\mu} v_{\nu}, \quad b_{\mu \nu}=\frac{1}{(\mu+\nu)^{1+\delta}} \quad(\delta>0) .
$$

Since $\Sigma \Sigma(\mu+\nu)^{-2-2 \delta}$ is convergent, it follows that

$$
\Sigma \Sigma \frac{u_{\mu} v_{\nu}}{(\mu+\nu)^{1+\delta}}
$$

is convergent. This is an imperfect form of a theorem to be proved later (Theorem 315).

The theorems concerning $\mathfrak{M}_{r}$ deduced from Hölder's inequality (Theorems 16 and 17) are unchanged, except that the statement of the second exceptional case of Theorem 16 must be modified in accordance with our remarks at the end of $\S 5.2$. Here we must say 'unless $s \leqq 0$ and $\mathfrak{M}_{s}(a)=0$ '.

One new point of greater interest arises in connection with this group of theorems. There is a theorem, suggested by Theorem 15 but not a corollary of it (even when it has been extended to infinite series), which has no analogue for finite sums.
161. ${ }^{\text {a }}$ If $k>1$ and $\Sigma a b$ is convergent for all $b$ for which $\Sigma b^{k^{\prime}}$ is convergent, then $\Sigma a^{k}$ is convergent.

We deduce this from another theorem due to Abel ${ }^{\mathrm{b}}$, which is of great interest in itself.
162. If $\Sigma a_{n}$ is divergent and

$$
A_{n}=a_{1}+a_{2}+\ldots+a_{n},
$$

${ }^{a}$ Landau (1).
${ }^{\mathrm{b}}$ Abel (1). There are theorems of the same type involving an arbitrary function $f(x)$. Thus, if $\Sigma a_{n}$ is divergent, $f(x)$ is positive and decreasing, and

$$
I=\int_{1}^{\infty} f(x) d x,
$$

then the convergence of $I$ involves that of $\Sigma a_{n} f\left(A_{n}\right)$, and the divergence of $I$ involves that of $\sum a_{n} f\left(A_{n-1}\right)$ : see, for example, de la Vallée Poussin (1, 398-399), Littlewood (1). This theorem, though of a more general character, does not actually include Theorem 162: it is not true that the divergence of $I$ necessarily involves that of $\Sigma a_{n} f\left(A_{n}\right)$. For an example to the contrary take

$$
a_{n}=2^{2^{n}}, f(x)=\frac{1}{x \log x} .
$$

so that $A_{n} \rightarrow \infty$, then
(i) $\Sigma \frac{a_{n}}{A_{n}}$ is divergent,
(ii) $\Sigma \frac{a_{n}}{A_{n}{ }^{1+\delta}}$ is convergent for every positive $\delta$.
(i) We have

$$
\frac{a_{n+1}}{A_{n+1}}+\frac{a_{n+2}}{A_{n+2}}+\ldots+\frac{a_{n+r}}{A_{n+r}} \geqq \frac{A_{n+r}-A_{n}}{A_{n+r}}=1-\frac{A_{n}}{A_{n+r}},
$$

which tends to 1 when $n$ is fixed and $r \rightarrow \infty$, and is therefore greater than $\frac{1}{2}$ for any $n$ and some corresponding $r$. This proves (i).
(ii) We may obviously suppose $0<\delta<1$. Then

$$
\Sigma \frac{A_{n}{ }^{\delta}-A_{n-1}{ }^{\delta}}{A_{n-1}{ }^{\delta} A_{n}{ }^{\delta}}=\Sigma\left(\frac{1}{A_{n-1}{ }^{\delta}}-\frac{1}{A_{n}{ }^{\delta}}\right)
$$

is convergent. By Theorem 41, the numerator on the left is not less than $\delta A_{n}{ }^{\delta-1}\left(A_{n}-A_{n-1}\right)=\delta a_{n} A_{n}{ }^{\delta-1}$. It follows that

$$
\Sigma \frac{a_{n}}{A_{n} A_{n-1}^{\delta}}
$$

is convergent. We prove in fact a little more than (ii).
To deduce Theorem 161, write

$$
a^{k}=u, \quad a b=u v, \quad b^{k^{\prime}}=u v^{k^{\prime}} .
$$

We then have to prove that, if $\Sigma u_{n}$ is divergent, there is a $v_{n}$ such that $\Sigma u_{n} v_{n}$ is divergent and $\Sigma u_{n} v_{n}{ }^{k^{\prime}}$ convergent. We take $v_{n}=1 / U_{n}$, where $U_{n}=u_{1}+u_{2}+\ldots+u_{n}$, and the conclusion follows from Theorem 162.
5.5. The means $\mathfrak{M}_{r}$ (continued). There is little to be added about the means $\mathfrak{M}_{r}$, but one or two further remarks are required. We begin with a remark concerning the generalisation of Theorem 4. This theorem, in so far as it concerns positive $r$, must be interpreted as follows: if the $a$ are bounded, and $a^{*}=\operatorname{Max} a$ is their upper bound, then

$$
\mathfrak{M}_{r} \rightarrow \mathfrak{M}_{\infty}=a^{*},
$$

when $r \rightarrow+\infty$; if the $a$ are unbounded, but $\mathfrak{M}_{r}$ is finite for all positive $r$, then $\mathfrak{M}_{r} \rightarrow \infty$.

The question of the continuity of $\mathfrak{M}_{r}$ for a finite positive or negative $r$ is no longer quite trivial. We state a comprehensive theorem, but give no proof, since all the points involved arise in a more interesting form in Ch. VI (§§6.10-6.11).

If $a_{n}=C, \mathfrak{M}_{r}=C$ for all $r$ (whether $C$ be positive or zero). We exclude this case, and also the case in which $\mathfrak{F}$ is meaningless, when $\mathfrak{M}_{r}=\infty$ for $r>0$ and $\mathfrak{M}_{r}=0$ for $r<0$. We write

$$
\mathfrak{Z}_{r}(a)=\log \mathfrak{M}_{r}(a)
$$

(with the conventions $\log \infty=+\infty, \log 0=-\infty$ ).
163. Apart from the cases just mentioned, the set of values I for which $\Omega_{r}(a)$ is finite is either the nul set or a closed, half-closed, or open interval $(u, v)$, where $-\infty \leqq u \leqq v \leqq+\infty$, which has $r=0$ as an internal or end-point, so that $u \leqq 0 \leqq v$, but is otherwise arbitrary (and in particular may include all real values or none). $\mathfrak{L}_{r}$ is $+\infty$ to the right, and $-\infty$ to the left, of $I$; is a continuous and strictly increasing function of $r$ in the interior (if it exists) of $I$; and tends, when $r$ approaches an end-point of I from inside I, to a limit equal to its value at the end-point.
5.6. The sums $\Im_{r}$. The definition of $\Im_{r}$ given in $\S 2.10$ is unchanged, and there is little to be said about the theorems concerning it, though those which involve continuity in $r$ are naturally less immediate. Theorem 20 must beinterpreted as meaning 'if $\mathbb{S}_{r}$ is convergent for some (sufficiently large) $r$, then it is convergent for all greater $r$ and ...', and Theorem 21 as meaning 'if $\widetilde{S}_{r}$ is convergent for all positive $r$ (however small) then...'. The extension of Theorem 20 may be proved as follows. If $\mathfrak{S}_{R}$ is convergent for a positive $R$, then $a_{n} \rightarrow 0$, and $\mathfrak{S}_{r}$ is convergent for $r>R$. There is a largest $a$, which we may suppose on grounds of homogeneity to be 1 , and we may suppose the $a$ arranged in descending order. If then

$$
\begin{gathered}
a_{1}=a_{2}=\ldots=a_{N}=1>a_{N+1} \\
\mathfrak{\Im}_{r}=\left(N+a_{N+1}^{r}+a_{N+2}^{r}+\ldots\right)^{1 / r}
\end{gathered}
$$

we have
for $r>R$. The series here lies between 1 and.

$$
N+a_{N+1}^{R}+a_{N+2}^{R}+\ldots
$$

from which the theorem follows.
There is one new theorem (trivial in the finite case).
164. If $\Im_{R}$ is convergent, then $\Im_{r}$ is continuous for $r>R$ and continuous on the right for $r=R$. If $\Im_{R}$ is divergent, but $\Im_{r}$ convergent for $r>R$, then $\mathfrak{S}_{r} \rightarrow \infty$ when $r \rightarrow R$.

The proof may be left to the reader.
5.7. Minkowski's inequality. The main arguments of $\S \S 2.11-2.12$ require no alteration. Theorems 24-26 suggest a further generalisation, with both summations infinite.
165. ${ }^{\text {a }}$ If $r>1$ and $a_{m n}$ is not of the form $b_{m} c_{n}$, then

$$
\left\{\sum_{n} q_{n}\left(\sum_{m} p_{m} a_{m n}\right)^{r}\right\}^{1 / r}<\sum_{m} p_{m}\left(\underset{n}{ } \sum_{n} q_{m n}^{r}\right)^{1 / r}
$$

The inequality is reversed when $0<r<1$.
There is no real loss of generality in supposing $p=1, q=1$, and the proof goes as before. Similarly, corresponding to Theorem 27, we have
166.

$$
\sum_{n}\left(\sum_{m} a_{m n}\right)^{r}>\sum_{m} \sum_{n} a_{m n}^{r}
$$

if $r>1$ (with a reversed inequality if $0<r<1$ ), unless, for every $n$, $a_{m n}=0$ for every $m$ save one.
5.8. Tchebychef's inequality. As one further illustration we take Tchebychef's inequality (Theorem 43).

We may suppose $\Sigma p=1$. The identity

$$
\sum_{1}^{n} p_{\mu} \sum_{1}^{n} p_{\nu} a_{\nu} b_{\nu}-\sum_{1}^{n} p_{\mu} a_{\mu} \sum_{1,}^{n} p_{\nu} b_{\nu}=\frac{1}{2} \sum_{1}^{n} \sum_{1}^{n} p_{\mu} p_{\nu}\left(a_{\mu}-a_{\nu}\right)\left(b_{\mu}-b_{\nu}\right)
$$

shows, provided of course that neither (a) nor (b) is nul, (i) that if $(a)$ and $(b)$ are similarly ordered then the convergence of $\Sigma p a b$ implies that of $\Sigma p a$ and $\Sigma p b$, and (ii) that if (a) and (b) are oppositely ordered then the convergence of $\Sigma p a$ and $\Sigma p b$ implies that of $\Sigma p a b$. In either case we may put $n=\infty$ in the identity, and our conclusions follow as before.
5.9. A summary. The theorem which follows is substantially an enumeration of the principal theorems of Ch. II which remain valid, with the glosses which we have explained in the preceding sections, for infinite series.

[^33]167. Theorems 1, 2, 3, 4, 5, 7, 9, 10, 11, 12, 13, 14, 15, 16, $17,18,19,20,21,22,24,25,27,28$ and 43 remain valid when the series concerned are infinite, provided that the inequalities asserted are interpreted in accordance with the conventions laid down in §5.1, and that the statement of the exceptional cases in Theorems 1, 2, 5, 10, 16, 24 and 25 is modified in the manner explained in §5.2.

It may be worth while to supplement the last clause of the theorem by a more explicit statement. The last words of the theorems must be replaced by

$$
\text { (1) 'or else } r<0 \text { and } \mathfrak{M}_{r}(a)=0 \text { ', }
$$

(2) 'or $\mathfrak{G}(a)=0$ ',
(5) 'or $r \leqq 0$ and $\mathfrak{M}_{r}(a)=0$ ',
(10) 'or (2) $\mathfrak{G}(a+b+\ldots+l)=0$ ',
(16) 'or $s \leqq 0$ and $\mathfrak{M}_{s}(a)=0$ ',
(24) 'or $r \leqq 0$ and $\mathfrak{M}_{r}(a+b+\ldots+l)=0$ ',
(25) 'or $r<0$ and $\left(\Sigma(a+b+\ldots+l)^{r}\right)^{1 / r}=0$ '.

We may add also that (as is explained in §6.4) most of the theorems referred to in Theorem 167 (especially those concerning $\mathfrak{M}_{r}$ ) may be derived by specialisation from the corresponding theorems for integrals. In Ch. VI, however, we often ignore negative values of $r$.

## MISCELLANEOUS THEOREMS AND EXAMPLES

The theorems which follow are for the most part connected with Theorems 156 and 157. We suppose in Theorems 168-175 that $f(x)$ and $g(x)$ are inverse functions which vanish for $x=0$ and increase steadily as $x$ increases, and that

$$
F(x)=\int_{0}^{x} f(u) d u, \quad G(x)=\int_{0}^{x} g(t) d t .
$$

168. If $\Sigma F\left(a_{n}\right)$ and $\Sigma G\left(b_{n}\right)$ are convergent, then $\Sigma a_{n} b_{n}$ is convergent, and

$$
\Sigma a_{n} b_{n} \leqq \Sigma F^{\prime}\left(a_{n}\right)+\Sigma G\left(b_{n}\right)
$$

[Corollary of Theorem 156.]
169. If $\Sigma a_{n} f\left(a_{n}\right)$ and $\Sigma b_{n} g\left(b_{n}\right)$ are convergent, then $\Sigma a_{n} b_{n}$ is convergent, and $\quad \Sigma a_{n} b_{n} \leqq \Sigma a_{n} f\left(a_{n}\right)+\Sigma b_{n} g\left(b_{n}\right)$.
[Corollary of Theorem 157.]
170. It is possible to choose $f$ (and so $g, F, G$ ) and $a_{n}$ in such a manner that $\Sigma F\left(a_{n}\right)$ is divergent, but $\Sigma a_{n} b_{n}$ is convergent for all $b_{n}$ for which $\Sigma G\left(b_{n}\right)$ is convergent.
171. It is also possible to make $\Sigma a_{n} f\left(a_{n}\right)$ divergent, but $\Sigma a_{n} b_{n}$ convergent whenever $\Sigma b_{n} g\left(b_{n}\right)$ is convergent.
[The point of the last two theorems is to show that Theorems 168 and 169 have no converses in the sense in which Theorem 161 gives a
converse of Hölder's inequality and the convergence test deduced fromit. Theorem 171 is proved by Cooper (3), and Theorem 170, which includes Theorem 171 and is a little stronger, may be proved in the same manner.]
172. If $\Sigma \frac{b_{n}}{\log \left(1 / b_{n}\right)}$ is convergent, then $\Sigma \frac{b_{n}}{\log n}$ is convergent.
[Cooper (3): Theorem 172 is used in Cooper's proof of Theorem 171.]
173. If $g(x)$ satisfies the inequality

$$
g(x y) \leqq g(x) g(y),
$$

and if $\Sigma a_{n} b_{n}$ is convergent whenever $\Sigma b_{n} g\left(b_{n}\right)$ is convergent, then $\Sigma a_{n} f\left(a_{n}\right)$ is convergent. Similarly, if $\Sigma a_{n} b_{n}$ is convergent whenever $\Sigma G\left(b_{n}\right)$ is convergent, then $\Sigma F\left(a_{n}\right)$ is convergent.
[See Cooper (3) for the first form, which in this case is stronger, the second form being a corollary.]
174. If $\Sigma a_{n} b_{n}$ is convergent whenever $\Sigma G\left(b_{n}\right)$ is convergent, then there is a number $\lambda=\lambda(a)$, depending upon the sequence ( $a$ ), for which $\Sigma F\left(\lambda a_{n}\right)$ is convergent.
175. If the conditions of Theorem 174 are satisfied, and $F^{\prime}(c x) \leqq k F^{\prime}(x)$ for small $x$, a $c>1$, and some $k$, then $\Sigma F\left(a_{n}\right)$ is convergent.
[For the last two theorems see Birnbaum and Orlicz (1).]
176. If $a_{n}$ and $b_{n}$ tend to zero, $k$ is positive, and

$$
\Sigma \frac{a_{n}}{\log \left(1 / a_{n}\right)}, \quad \Sigma e^{-k / b_{n}}
$$

are convergent, then $\Sigma a_{n} b_{n}$ is convergent.
[Use Theorem 169.]
177. If $x>0, a_{n} \geqq 0$, and $f(x)=\Sigma a_{n} x^{n}$, then $f(x)$ is a convex function of $x$ and $\log f(x)$ of $\log x$.
[Plainly $f^{\prime \prime}(x) \geqq 0$. To prove the second result, let $x=e^{-y}, f(x)=g(y)$. Then

$$
g g^{\prime \prime}-g^{\prime 2}=\Sigma a_{n} e^{-n y} \Sigma n^{2} a_{n} e^{-n y}-\left(\Sigma n a_{n} e^{-n y}\right)^{2} \geqq 0,
$$

by Theorem 7. The result follows from Theorem 118.]
178. If $a_{n} \geqq 0, \lambda_{n}>\lambda_{n-1} \geqq 0$, and $f(x)=\Sigma a_{n} e^{-\lambda_{n} x}$, then $\log f(x)$ is a convex function of $x$.
179. If $a_{n}>0$ and $\lambda_{n}, \mu_{n}, \ldots, \nu_{n}, x, y, \ldots, z$ are real, then the domain $D$ of convergence of the series

$$
\Sigma a e^{\lambda x+\mu y+\ldots+\nu z}=f(x, y, \ldots, z)
$$

is convex, and $\log f$ is a convex function of $x, y, \ldots, z$ in $D$.
[Because (by the extension of Theorem 11 to infinite series)

$$
f\left\{x_{1} t+x_{2}(1-t), \ldots, z_{1} t+z_{2}(1-t)\right\} \leqq\left\{f\left(x_{1}, \ldots, z_{1}\right)\right\}^{t}\left\{f\left(x_{2}, \ldots, z_{2}\right)\right\}^{1-t} .
$$

Here our conventions concerning convergence are important.]
180.

$$
\Sigma a_{n}^{2}<2\left(\Sigma n^{2} a_{n}^{2}\right)^{\frac{1}{2}}\left(\Sigma\left(a_{n}-a_{n+1}\right)^{2}\right)^{\frac{1}{2}},
$$

unless $a_{n}=0$ for all $n$.
[See Theorem 226.]

## CHAPTER VI

## INTEGRALS

6.1. Preliminary remarks on Lebesgue integrals. The integrals considered in this chapter are Lebesgue integrals, except in $\S \S 6.15-6.22$, where we are concerned with Stieltjes integrals. It may be convenient that we should state here how much knowledge of the theory we assume. This is for the most part very little, and all that the reader usually needs to know is that there is some definition of an integral which possesses the properties specified below. There are naturally many of our theorems which remain significant and true with the older definitions, but the subject becomes easier, as well as more comprehensive, if the definitions presupposed have the proper degree of generality.

We take for granted the idea of a measurable set, usually in one but occasionally in more dimensions. The sets which we consider may be bounded or unbounded. The definition of measure applies in the first instance to bounded sets: an unbounded set is said to be measurable if any bounded part of it is measurable, and its measure is the upper bound of the measures of its bounded components.

We shall generally assume, without special remark, that any set $E$ with which we are concerned is measurable. We denote the measure of $E$ by $m E$ or sometimes, where there is no risk of ambiguity, simply by $E$. If $E$ is unbounded, $m E$ may be $\infty$.

We also take for granted the idea of a measurable function. Sums, products, and limits of measurable functions are measurable. All functions definable by the ordinary processes of analysis are measurable, and we shall confine our attention to measurable functions; we shall not usually repeat explicitly that a function which occurs in our work is assumed to be measurable.

Next, we take for granted the definition of the integral, of a bounded or unbounded function, over any (bounded or unbounded) interval or measurable set of points. A bounded measur-
able function is integrable over any bounded measurable set. We call the class of (bounded or unbounded) functions integrable over the interval or set $E$ in question the class $L$ or, if it is desirable to emphasise the set in question, the class $L(E)$. If $f$ belongs to $L$, we say that ' $f$ is $L$ ', and write $f \in L$ (and similarly for other classes).

If $f=1$ then

$$
\int_{E} f d x=m E .
$$

If $f \in L$, then $|f| \epsilon L$. If $f^{+}$and $f^{-}$are the functions equal to $f$ when $f$ is positive and negative respectively, and to zero otherwise, so that

$$
f^{+}=\operatorname{Max}(f, 0), \quad f^{-}=\operatorname{Min}(f, 0), \quad f=f^{+}+f^{-}, \quad|f|=f^{+}-f^{-}
$$

then $f^{+} \epsilon L$ and $f^{-} \epsilon L$, and ${ }^{\text {a }}$

$$
\int f d x=\int f^{+} d x+\int f^{-} d x, \quad \int|f| d x=\int f^{+} d x-\int f-d x
$$

If $f \geqq 0$, and $(f)_{n}=\operatorname{Min}(f, n)$, then

$$
\int f d x=\lim _{n \rightarrow \infty} \int(f)_{n} d x
$$

(substantially by definition).
If $f \in L$, and ( $g$ is measurable and) $|g|<C|f|$, then $g \in L$.
If $f_{1}, f_{2}, \ldots, f_{n} \in L$ then
$\int\left(a_{1} f_{1}+a_{2} f_{2}+\ldots+a_{n} f_{n}\right) d x=a_{1} \int f_{1} d x+a_{2} \int f_{2} d x+\ldots+a_{n} \int f_{n} d x$.
If $p>0$ and ( $f$ is measurable and) $|f|^{p} \in L$, we say that $f$ belongs to the Lebesgue class $L^{p}$, or $f \in L^{p}$. These classes are most important when $p \geqq 1 . L^{1}$ is identical with $L$.

If the integration is over a finite interval (or bounded set), then $L^{p}$ includes every $L^{q}$ for which $q>p ; f \in L^{q}$ implies $f \in L^{p}$. A bounded function belongs to every $L^{q}$. These propositions are not true for an infinite interval; $f$ may belong to $L^{p}$, in $(0, \infty)$, for one value of $p$ only.

If the interval is finite and $f \in L^{q}, p<q$, then $|f|^{p}<1+|f|^{q}$, so that $f \in L^{p}$.

If the interval is $(0, a)$, where $a<1$, then ( $a$ ) $x^{-1 / p}$ belongs to $L^{p-\delta}$ for every $\delta>0$, but not to $L^{p}$; (b) $x^{-1 / p}\left(\log \frac{1}{x}\right)^{-2 / p}$ belongs to $L^{p}$, but not to

[^34]$L^{p+\delta} ;(c) \log (1 / x)$ belongs to every $L^{p}$; and (d) $e^{1 / x}$ belongs to no $L^{p}$. If the interval is $(0, \infty)$, then $x^{-\frac{1}{2}}(1+|\log x|)^{-1}$ belongs to $L^{2}$ but to no other class $L^{p}$.
6.2. Remarks on nul sets and nul functions. A set of measure zero is called a nul set. Nul sets are negligible in the theory of integration. If $f=g$ except in a nul set, we say that $f$ and $g$ are equivalent, and write $f \equiv g$. Equivalent functions have the same integral (if any).
If $f \equiv 0$, we call $f$ a nul function, and say that $f$ is nul.
Similarly, we define 'equivalent in $E$ ', 'nul in $E$ ': $f$ is nul in $E$ if $f=0$ at all points of $E$ except the points of a nul set. In such cases we shall not repeat the reference to $E$ when the context makes it obvious, as for example when we are considering integrals over $E$.

If a property $P(x)$ is possessed by all $x$ except the $x$ of a nul set, we shall say that it is possessed by almost all $x$, or that $P(x)$ is true for almost all $x$, or almost always. Thus a nul function is almost always zero.

We shall generally assume that our functions $f, g, \ldots$ are almost always finite; but there will be occasions when we have to consider functions infinite in a set of positive measure. Thus if $f$ is generally positive, but zero in a set $E$ of positive measure, and $r<0$, then we must regard $f^{r}$ as infinite in $E$, and $\int f^{r} d x$ as having the value $\infty$.

If $E$ is nul then

$$
\int_{E} f d x=0
$$

for all $f$. We shall assume without special remark that a set $E$ over which an integral is extended is not nul.

If $f \geqq 0$, then a necessary and sufficient condition that $\int f d x=0$ is that $f$ should be nul.

It may be worth while to call attention explicitly to the theorem which replaces this in the theory of Riemann integration. We denote the class of Riemann integrable functions by $R$. A necessary and sufficient condition that $f$ should be $R$ is that $f$ should be bounded and that its set of discontinuities should be nul.

If $f$ is $R$, and $f \geqq 0$, then a necessary and sufficient condition that $\int f d x=0$ is that $f=0$ at all points of continuity of $f$.

For (1) if the condition is satisfied, then $f \equiv 0$ and so $\int f d x=0$. And (2) if it is not satisfied, then there is a point of continuity $\xi$ at which $f(\xi)>0$,
and an interval including $\xi$ throughout which $f(x)>\frac{1}{2} f(\xi)$; so that $\int f d x>0$.

This theorem enables us to specify the cases of equality in our inequalities when they are restricted to functions of $R$. In fact most of our theorems have a dual interpretation. The primary interpretation is that in which the integrals are Lebesgue integrals and 'nul' and 'equivalent' are interpreted as in the theory of Lebesgue. In the secondary interpretation integrals are 'Riemann integrals', a 'nul function' is a function which is zero at all its points of continuity, and 'equivalent functions' are functions whose difference is nul in this sense.
6.3. Further remarks concerning integration. What has beensetoutin $\S \S 6.1$ and 6.2 is a sufficientfoundationfor most of the subsequent theorems: for example, for the most complete forms of Hölder's and Minkowski's inequalities (Theorems 188 and 198). There will be a few occasions on which we shall have to appeal to more difficult theorems, and we enumerate these here.
(a) Integration by parts. The theorem required is: iff and $g$ are integrals (absolutely continuous functions), then

$$
\int_{a}^{b} f g^{\prime} d x=[f g]_{a}^{b}-\int_{a}^{b} f^{\prime} g d x .
$$

(b) Passage to the limit under the integral sign. The two main theorems are
(i) If $\left|s_{n}(x)\right|<\phi(x)$, where $\phi \in L$, and $s_{n}(x)$ tends to a limit s $(x)$ for all or almost all $x$, then

$$
\int s_{n}(x) d x \rightarrow \int s(x) d x .
$$

(ii) If $s_{n}(x) \in L$ for every $n, s_{n}(x)$ increases with $n$ for all or almost all $x$, and

$$
\lim s_{n}(x)=s(x)
$$

then

$$
\int s_{n}(x) d x \rightarrow \int s(x) d x
$$

In (ii) the integral on the right may be infinite, when the result is to be interpreted as $\int s_{n}(x) d x \rightarrow \infty$; in particular this happens if $s(x)=\infty$ in a set of positive measure. In each of these theorems $n$ may be an integer which tends to infinity or a continuous parameter which tends to a limit.

It follows from (i), as is shown in books on the theory of functions of a real variable, that a function $f(x)$ whose incrementary ratio

$$
\frac{f(x+h)-f(x)}{h}
$$

is bounded (and which therefore has a derivative almost everywhere) is the integral of its derivative. Combining this remark with that at the end of § 3.18 we see that a continuous convex function $f(x)$ is the integral of its derivative $f^{\prime}(x)$, or of its one-sided derivatives $f_{l}^{\prime}(x), f_{r}^{\prime}(x)$. It is therefore the integral of an increasing function. On the other hand, if $f(x)$ is the integral of an increasing function $g(x)$, and $h>0$, then

$$
f(x+h)-f(x)=\int_{x}^{x+h} g(u) d u \geqq \int_{x-h}^{x} g(u) d u=f(x)-f(x-h),
$$

so that $f(x)$ is convex. Hence the class of continuous convex functions is identical with that of integrals of increasing functions.
An increasing function belongs to $R$, so that the integrals in question exist as Riemann integrals, and the theorem could be proved without any use of the theory of Lebesgue.
(c) Substitution. The standard theorem is: if $f$ and $g$ are integrable, $g \geqq 0, G$ is an integral of $g$, and $a=G(\alpha), b=G(\beta)$, then

$$
\begin{equation*}
\int_{a}^{b} f(x) d x=\int_{a}^{\beta} f\{G(y)\} g(y) d y \tag{6.3.1}
\end{equation*}
$$

Here any of $a, b, \alpha, \beta$ may be infinite.
This theorem covers all cases in which we shall require the rule for transformation of an integral by change of the independent variable ${ }^{\text {a }}$. But we shall generally use only trivial substitutions such as $x=y+a$ or $x=a y$, when the validity of the rule follows at once from the definitions.
(d) Multiple and repeated integrals. The only theorem appealed to is 'Fubini's Theorem'. If $f(x, y)$ is (measurable and) nonnegative, and any one of the integrals

$$
\int_{a}^{A} \int_{b}^{B} f d x d y, \quad \int_{a}^{A} d x \int_{b}^{B} f d y, \quad \int_{b}^{B} d y \int_{a}^{A} f d x
$$

exists, then the other integrals exist, and all are equal. Here the limits are finite or infinite, and the case of divergence is included; if one integral diverges the others diverge.
a We may add two additional remarks concerning the formula (6.3.1).
(1) If we suppose, as in the text, that $g$ is non-negative and integrable, but assume only the measurability (and not the integrability) of $f$, then the existence of the right-hand side of (6.3.1) is a sufficient, as well as a necessary, condition for the existence of the left-hand side, i.e. for the integrability of $f$.
(2) The integrability of $f(x)$, though it implies that of $f\{G(y)\} g(y)$, does not imply even the measurability of $f\{G(y)\}$.

Suppose then that $f(x, y)$ is measurable and non-negative. The double integral is zero if and only if $f(x, y)$ is nul, i.e. if the set in which $f(x, y)>0$ has measure zero. The first repeated integral is zero if and only if $f(x, y)$ is, for almost all $x$, nul in $y$; and the second if $f(x, y)$ is, for almost all $y$, nul in $x$. Hence these three senses of 'a nul non-negative function of two variables' are equivalent.
6.4. Remarks on methods of proof. Inequalities proved for finite sums may often be extended to integrals by the use of limiting processes, but something is usually lost in the argument. We may illustrate this by considering the analogue for integrals of Theorem 7.
Suppose first that $f(x)$ and $g(x)$ are non-negative and Riemann integrable in ( 0,1 ); and take

$$
a_{\nu}=f\left(\frac{\nu}{n}\right), \quad b_{\nu}=g\left(\frac{\nu}{n}\right)
$$

in Theorem 7. Dividing by $n^{2}$, we obtain ${ }^{\text {a }}$

$$
\left\{\frac{1}{n} \Sigma f\left(\frac{\nu}{n}\right) g\left(\frac{\nu}{n}\right)\right\}^{2} \leqq \frac{1}{n} \Sigma f^{2}\left(\frac{\nu}{n}\right) \cdot \frac{1}{n} \Sigma g^{2}\left(\frac{\nu}{n}\right) ;
$$

and, making $n \rightarrow \infty$,

$$
\begin{equation*}
\left(\int_{0}^{1} f g d x\right)^{2} \leqq \int_{0}^{1} f^{2} d x \int_{0}^{1} g^{2} d x \tag{6.4.1}
\end{equation*}
$$

If we use the Lebesgue integral we must argue differently ${ }^{b}$. Suppose that $f$ and $g$ are non-negative and $L^{2}$ in ( 0,1 ), and that $e_{r s}$ is the set in which

$$
\frac{r-1}{n} \leqq f^{2}<\frac{r}{n}, \quad \frac{s-1}{n} \leqq g^{2}<\frac{s}{n} \quad(r, s=1,2,3, \ldots) .
$$

Then

$$
\left(\int_{0}^{1} f g d x\right)^{2} \leqq\left[\Sigma \Sigma\left(\frac{r}{n}\right)^{\frac{1}{2}}\left(\frac{s}{n}\right)^{\frac{1}{e}} e_{r s}\right]^{2} \leqq \Sigma \Sigma \frac{r}{n} e_{r s} \Sigma \Sigma \frac{s}{n} e_{r s},
$$

by Theorem 7. Now

$$
\Sigma \Sigma \frac{r}{n} e_{r s}=\Sigma \Sigma{ }_{n}^{1} e_{r s}+\Sigma \Sigma \frac{r-1}{n} e_{r s} \leqq \frac{1}{n}+\int_{0}^{1} f^{2} d x
$$

[^35]and there is a similar inequality involving $g$. Hence, making $n \rightarrow \infty$, we obtain (6.4.1).

In either case our final result is imperfect. Even if we can use Theorem 7 with ' < ', this will degenerate into ' $\leqq$ ' when we pass to the limit, and we shall lose touch with the cases of equality.

The passage in the opposite direction, from an integral inequality to an inequality for sums, is much simpler, and can be effected by suitable specialisation. Consider, for example, the inequality

$$
\begin{equation*}
\exp \left\{\int_{0}^{1} \log f(x) d x\right\}<\int_{0}^{1} f(x) d x \tag{6.4.2}
\end{equation*}
$$

(§6.7, Theorem 184). If

$$
q_{\nu}>0, \quad q_{1}+q_{2}+\ldots+q_{n}=1
$$

and we define $f(x)$ by

$$
f(x)=a_{\nu} \quad\left(q_{1}+\ldots+q_{\nu-1} \leqq x<q_{1}+\ldots+q_{\nu-1}+q_{\nu}\right)
$$

it being understood that $q_{1}+\ldots+q_{\nu-1}$ means 0 when $\nu=1$, we obtain Theorem 9. The conditions under which inequality degenerates into equality in Theorem 9 also follow immediately from the corresponding conditions for (6.4.2).

This method of proof is often useful, since integrals are often more manageable than series. We shall meet with examples in Ch. IX.
6.5. Further remarks on method: the inequality of Schwarz. We meet the difficulty of $\S 6.4$, as in our treatment of infinite series, by going back to the proofs of the theorems of Ch. II, and observing that, with the obvious changes, they can be applied to integrals of the most general type. We mayillustrate the point here by considering 'Schwarz's' inequality', the analogue of Theorem 7.
181. $\left(\int f g d x\right)^{2}<\int f^{2} d x \int g^{2} d x$, unless $A f \equiv B g$, where $A$ and $B$ are constants, not both zero.

Here, and later, we suppress the limits of integration when there is nothing to be gained by showing them explicitly; they

[^36]may be finite or infinite, or the integrals may be over any measurable set $E$, in which case of course $A f \equiv B g$ means $A f \equiv B g$ in $E$. We also adopt conventions corresponding to those of §5.1: ' $X<Y$ ' means 'if $Y$ is finite then $X$ is finite and $X<Y$ '; and inequalities of other forms like those mentioned in §5.1 are to be interpreted similarly. Thus every inequality contains implicitly an assertion about 'convergence', which we shall only make explicit occasionally. For example, Theorem 181 asserts implicitly that 'if $\int f^{2} d x$ and $\int g^{2} d x$ are finite, then $\int f g d x$ is finite; if $f$ and $g$ are $L^{2}$, then $f g$ is $L^{\prime}$.

The proofs corresponding to those of $\S 2.4$ run as follows.
(i) We have

$$
\begin{aligned}
& \int f^{2} d x \int g^{2} d x-\left(\int f g d x\right)^{2} \\
& \quad=\frac{1}{2} \int f^{2}(x) d x \int g^{2}(y) d y+\frac{1}{2} \int g^{2}(x) d x \int f^{2}(y) d y \\
& \quad=\frac{1}{2} \int d y \int\{f(x) g(y)-g(x) f(y)\}^{2} d x \geqq 0 .
\end{aligned}
$$

It remains to discuss the possibility of equality. In the first place, there is certainly equality if $A f \equiv B g$. Next, if there is equality, and $g$ is nul, then $A f \equiv B g$ with $A=0, B=1$. We may therefore assume that $g$ is not nul, so that the set $E$ in which $g(y) \neq 0$ has positive measure. If

$$
\int d y \int\{f(x) g(y)-g(x) f(y)\}^{2} d x=0
$$

then

$$
\begin{equation*}
\int\{f(x) g(y)-g(x) f(y)\}^{2} d x=0 \tag{6.5.1}
\end{equation*}
$$

for almost all $y$, and therefore for some $y$ belonging to $E$. We may therefore suppose that $g\left(y_{0}\right) \neq 0$ and that (6.5.1) is true for $y=y_{0}$. But then $\quad f(x) g\left(y_{0}\right)-g(x) f\left(y_{0}\right)=0$
for almost all $x$, and this completes the proof.
(ii) The quadratic form

$$
\int(\lambda f+\mu g)^{2} d x=\lambda^{2} \int f^{2} d x+2 \lambda \mu \int f g d x+\mu^{2} \int g^{2} d x
$$

is positive. We can now complete the proof as in §2.4.
The analogue of Theorem 181 for multiple integrals may be proved similarly. We shall not usually mention such extensions explicitly, but we shall occasionally take them for granted. It
is to be understood that, when we do this, the extension may be proved in the same manner as the original theorem.

We can translate the proof of Theorem 8 in the same manner, and so obtain
182.a $\left|\begin{array}{cccc}\int f^{2} d x & \int f g d x & \ldots & \int f h d x \\ \ldots \ldots & \ldots \ldots & \ldots & \ldots \ldots \\ \int h f d x & \int h g d x & \ldots & \int h^{2} d x\end{array}\right|>0$,
unless the functions $f, g, \ldots, h$ are linearly dependent, i.e. unless there are constants $A, B, \ldots, C$, not all zero, such that

$$
A f+B g+\ldots+C h \equiv 0
$$

## MEANS $\mathfrak{M}_{r}(f)$

6.6. Definition of the means $\mathfrak{M}_{r}(f)$ when $r \neq 0$. In what follows the sign of integration, used without specification of the range, refers to a finite or infinite interval $(a, b)$ or to a measurable set $E .{ }^{\mathrm{b}} f(x)$ is finite almost everywhere in $E$ and non-negative; $p(x)$, the 'weight function', is finite and positivec everywhere in $E$, and integrable over $E$. The parameter $r$ is real and not zero.

Our hypotheses involve $0<\int p d x<\infty$. It is often convenient to suppose

$$
\int p d x=1:
$$

in this case (cf. §2.2) we write $q$ for $p$.
We write

$$
\begin{align*}
\mathfrak{M}_{r}(f)=\mathfrak{M}_{r}(f, p) & =\left(\frac{\int p f^{r} d x}{\int p d x}\right)^{1 / r} \quad(r \neq 0),  \tag{6.6.1}\\
\mathfrak{H}(f) & =\mathfrak{M}_{\mathbf{1}}(f) \tag{6.6.2}
\end{align*}
$$

so that

$$
\begin{equation*}
\mathfrak{M}_{r}(f)=\left\{\mathfrak{U}\left(f^{r}\right)\right\}^{1 / r} ; \tag{6.6.3}
\end{equation*}
$$

with the following conventions. If $\int p f^{r} d x$ is infinite, we write

$$
\int p f^{r} d x=\infty, \quad \mathfrak{M}_{r}(f)=\infty \quad(r>0), \quad \mathfrak{M}_{r}(f)=0 \quad(r<0) .
$$

a Gram (1).
${ }^{b}$ When $r>0$ we can reduce every case to that of the interval ( $-\infty, \infty$ ), by supposing $f=0$ in the set complementary to $E$.
c The hypothesis $p \geqq 0$, instead of $p>0$, would lead to slightly different results concerning the cases of equality (e.g. $p f \equiv p C$ instead of $f \equiv C$ ). This case could be reduced to the apparently more special case by replacing $E$ by the sub-set of $E$ in which $p>0$.

In particular $\mathfrak{M}_{r}(f)=0$ if $r<0$ and $f=0$ in a set of positive measure. If we agree further to regard 0 and $\infty$ as reciprocals of one another, we have

$$
\begin{equation*}
\mathfrak{M}_{r}(f)=\frac{1}{\mathfrak{M}_{-r}(\mathbf{l} / f)} . \tag{6.6.4}
\end{equation*}
$$

This formula enables us to pass from positive to negative $r$, and we shall simplify the following theorems by restricting ourselves, for the most part, to positive $r$.

If $f \equiv 0, \mathfrak{M}_{r}(f)=0$ for all $r$. If $f \equiv C$, where $C$ is positive and finite, then $\mathfrak{M}_{r}(f)=C$ for all $r$. If $f \equiv \infty, a$ anen $\mathfrak{M}_{r}(f)=\infty$ for all $r$. Apart from these cases, $\mathfrak{M}_{r}(f)=\infty$ is possible only when $r>0$, and $\mathfrak{M}_{r}(f)=0$ when $r<0$.

We define $\operatorname{Max} f$, the 'effective upper bound' of $f$, as the largest $\xi$ which has the property:
'if $\epsilon>0$, there is a set $e(\epsilon)$ of positive measure in which $f>\xi-\epsilon$ '. If there is no such $\xi$, we write $\operatorname{Max} f=\infty$. For functions continuous in a closed interval, $\operatorname{Max} f$ is the ordinary maximum. $\operatorname{Min} f$ is defined similarly; $\operatorname{Min} f \geqq 0$ and

$$
\begin{equation*}
\operatorname{Min} f=\frac{1}{\operatorname{Max}(1 / f)} . \tag{6.6.5}
\end{equation*}
$$

Equivalent functions have the same Max and Min.
Suppose for example that the range of integration is $(0, \infty)$, and that $f(x)$ and $q(x)$ are the step functions defined by

$$
f(x)=a_{n}, \quad q(x)=q_{n} \quad(n-1 \leqq x<n, \quad n=1,2,3, \ldots) .
$$

Then

$$
\mathfrak{M}_{r}(f)=\left(\Sigma q a^{r}\right)^{1 / r}=\mathfrak{M}_{r}(a),
$$

according to the definition of (5.2.1). Similarly

$$
\operatorname{Max} f=\operatorname{Max} a, \quad \operatorname{Min} f=\operatorname{Min} a
$$

and (if we anticipate the definition of §6.7) $\mathfrak{G}(f)=\mathfrak{G}(a)$. This specialisation enables us to include many theorems of Ch. II and Ch. V in the corresponding theorems of this chapter.

Alternatively we might (asin § 6.4) suppose that the range of integration is ( 0,1 ), and define $f(x)$ and $q(x)$ by

$$
f(x)=a_{n} \quad\left(q_{1}+\ldots+q_{n-1} \leqq x<q_{1}+\ldots+q_{n}\right), \quad q(x)=1
$$

In this case also $\mathfrak{M}_{r}(f)$ reduces to $\mathfrak{M}_{r}(a)$.
a. To admit this case is to abandon momentarily the understanding of $\S 6.2$, that $f$ is assumed to be finite almost everywhere.
183. If $r \neq 0$ and $\mathfrak{M e}_{r}(f)$ is finite and positive, then

$$
\operatorname{Min} f<\operatorname{Me}_{r}(f)<\operatorname{Max} f
$$

unless $f \equiv C$.
Here $r$ is of either sign. The proof is like that of Theorem 1. Suppose first that $r=1$. Then, using a weight function $q(x)$, we have

$$
\int q(f-\mathfrak{H}) d x=0
$$

Hence either $f \equiv \mathfrak{H}$ or $f-\mathfrak{A}$ is positive and negative, each in a set of positive measure. This proves the result for $r=1$, and we extend it to the general case by use of (6.6.3).

If we wish to state Theorem 183 in a form corresponding more exactly to that of Theorem 1 and its extension in $\mathrm{Ch} . \mathrm{V}^{\text {a }}$, we must say ' $\operatorname{Min} f<\mathfrak{M}_{r}(f)<\operatorname{Max} f$ unless $f \equiv C$ or else $r<0$ and $\mathfrak{M}_{r}(f)=0$ '. We have then two cases of equality corresponding exactly to those distinguished in $\S 5.2$, the 'primary' case in which $f \equiv C$, in which both inequalities reduce to equalities, and the 'secondary' case, occurring only for $r<0$, in which one inequality only reduces to an equality. This distinction recurs in many of our theorems, when $r<0$, as it recurred in Chs. II and V ; but it is less conspicuous here because we often ignore negative values of $r$.
6.7. The geometric mean of a function. We define the geometric mean $\mathfrak{G s}(f)$ by

$$
\begin{equation*}
\mathscr{S}(f)=\mathscr{G}(f, p)=\exp \left(\frac{\int p \log f d x}{\int p d x}\right), \tag{6.7.1}
\end{equation*}
$$

or

$$
\begin{equation*}
\log \mathfrak{G}(f)=\mathfrak{A}(\log f), \tag{6.7.2}
\end{equation*}
$$

so that, in particular, if $p=q, \int q d x=1$, we have

$$
\begin{equation*}
\mathfrak{J}(f)=\log \mathfrak{G}(f)=\int q \log f d x . \tag{6.7.3}
\end{equation*}
$$

Certain preliminary explanations are necessary.
Since $\log f$ is not necessarily positive, the possibilities concerning the convergence of $\mathfrak{F}$ are more complex than those which we have considered hitherto.

[^37]If we denote by $\mathfrak{S}^{+}$and $\mathfrak{S}^{-}$the integrals formed with $\log ^{+} f$ and $\log ^{-} f$, as $\mathfrak{F}$ is formed with $\log f$, a then there are four possibilities: (a) $\mathfrak{S}^{+}$and $\mathfrak{S}^{-}$both finite, (b) $\mathfrak{F}^{+}$finite, $\mathfrak{S}^{-}=-\infty$, (c) $\mathfrak{J}^{+}=\infty$, $\mathfrak{J}^{-}$finite, (d) $\mathfrak{I}^{+}=\infty, \mathfrak{J}^{-}=-\infty$. The four cases are exemplified by the functions

$$
x, \quad e^{-1 / x}, \quad e^{1 / x}, \quad \exp \left(\frac{1}{x^{2}} \sin \frac{1}{x}\right)
$$

in $(0,1)$, with $q(x)=1$.
If $\mathfrak{M}_{r}(f)$ is finite for some $r>0$, then, since

$$
\log ^{+} f \leqq \operatorname{Max}\left(\frac{f^{r}-1}{r}, 0\right),
$$

$\mathfrak{S}^{+}(f)$ will be finite, and we shall be concerned only with cases $(a)$ and $(b)$. In case $(a), \mathfrak{F}(f)$ exists as a Lebesgue integral, and $\mathscr{F}(f)$ is positive and finite. In case (b) we write

$$
\mathfrak{F}(f)=-\infty, \quad \mathfrak{G}(f)=0 .
$$

Similarly, if $\mathfrak{M}_{r}(1 / f)$ is finite for some $r>0$, we are in case (a) or case (c); in the latter we write

$$
\mathfrak{J}(f)=\infty, \quad \mathscr{S}(f)=\infty .
$$

In case $(d)$ the symbol $\mathscr{G}(f)$ is meaningless. In this case $\mathfrak{M}_{r}(f)$ and $\mathfrak{M}_{r}(\mathbf{1} / f)$ are infinite for every positive $r$, and $\mathfrak{M}_{r}(f)=0$ for every negative $r$.

In case (a) we have

$$
\begin{equation*}
\mathscr{S}\left(\frac{1}{f}\right)=\frac{1}{\mathfrak{G}(f)}, \tag{6.7.4}
\end{equation*}
$$

both sides being positive and finite; and a moment's consideration shows that this equation holds in all cases, if we adopt the same convention about 0 and $\infty$ as in (6.6.4), and the additional convention that one side of the equation is meaningless if the other is meaningless.

We now prove the analogue of Theorem 9.
184. If $\mathfrak{A}(f)$ is finite then

$$
\begin{equation*}
\mathfrak{F}(f)<\mathfrak{A}(f), \tag{6.7.5}
\end{equation*}
$$

a See $\S 5.2$ for the definitions of $\log ^{+}$and $\log ^{-}$.
unless $f \equiv C$, where $C$ is constant. More generally, if $\mathfrak{M}_{r}(f)$, where $r>0$, is finite, then

$$
\begin{equation*}
\mathfrak{G}(f)<\mathfrak{M}_{r}(f), \tag{6.7.6}
\end{equation*}
$$

with the same reservation ${ }^{\text {a }}$.
Suppose first that $r=1, \mathfrak{M}_{r}=\mathfrak{A}$. If $\mathfrak{A}(f)=0, f \equiv 0$, and so $\mathfrak{J}(f)=-\infty, \mathfrak{G}(f)=0=\mathfrak{A}(f)$. We may thereforesuppose $\mathfrak{A}(f)>0$.

Since, by Theorem 150,

$$
\begin{equation*}
\log t<t-1 \tag{6.7.7}
\end{equation*}
$$

if $t>0, t \neq 1$, we have

$$
\begin{gathered}
\log f-\log \mathfrak{A}(f) \leqq \frac{f}{\mathfrak{A}(f)}-1, \\
\mathfrak{A}(\log f)-\log \mathfrak{A}(f) \leqq \frac{\mathfrak{A}(f)}{\mathfrak{A}(f)}-1=0, \\
\log \mathfrak{G}(f)=\mathfrak{A}(\log f) \leqq \log \mathfrak{A}(f) .
\end{gathered}
$$

Equality can occur only if $f \equiv \mathfrak{A}(f)$.
The result for general $r$ now follows from (6.6.3).
In Theorem 184 we have stated the hypotheses 'if $\mathfrak{A}(f)$ is finite', 'if $\mathfrak{M}_{r}(f)$ is finite' explicitly. As we have explained in $\S 5.1$ and 6.5 , we shall often save space by omitting such hypotheses in accordance with our conventions. We shall also denote constants by $C, A, B, a, b, \ldots$ without explanation, when there is no danger of ambiguity. Two $C$ 's occurring in the same connection will not necessarily be the same.

We add two corollaries (extensions of Theorem 10).
185. $\mathfrak{G}(f)+\mathfrak{G}(g)<\mathfrak{G}(f+g)$, unless $A f \equiv B g$, where $A, B$ are not both zero, or $\mathfrak{G}(f+g)=0$.

We may suppose that $\mathfrak{G}(f+g)>0$. Then, by Theorem 184,

$$
\frac{\mathfrak{G}(f)}{\mathfrak{G}(f+g)}=\mathfrak{G}\left(\frac{f}{f+g}\right) \leqq \mathfrak{A}\left(\frac{f}{f+g}\right) .
$$

The addition of the two inequalities of this type gives the result.
More generally
186. $\mathfrak{G}\left(f_{1}\right)+\mathfrak{G}\left(f_{2}\right)+\mathfrak{G}\left(f_{3}\right)+\ldots<\mathfrak{G}\left(f_{1}+f_{2}+f_{3}+\ldots\right)$ (the series being finite or infinite), unless $f_{n} \equiv C_{n} \Sigma f_{n}$ or $\mathfrak{G}\left(\Sigma f_{n}\right)=0$.
a For the proof which follows see F. Riesz (7).
6.8. Further properties of the geometric mean. Our next theorem corresponds to Theorem 3 (for positive $r$ ).
187. If $\mathfrak{M}_{r}(f)$ is finite for some positive $r$, then

$$
\begin{equation*}
\mathfrak{M}_{r}(f) \rightarrow \mathfrak{G}(f) \tag{6.8.1}
\end{equation*}
$$

when $r \rightarrow+0$.
It should be observed that $\mathbb{G}(f)$ may be finite even when $\mathfrak{M}_{r}(f)=\infty$ for all $r>0$. This is so, for example, if $f(x)=\exp \left(x^{-\frac{1}{2}}\right)$, $q(x)=1$, and the range is $(0,1)$.

When $E$ is a closed interval or set, and $f$ is continuous and positive, the proof is immediate. In this case $f \geqq \delta>0, \log f$ is bounded, and

$$
\begin{aligned}
\mathfrak{M}_{r}^{r}=\int e^{r \log f} q d x & =\int\left[1+r \log f+O\left\{r^{2}(\log f)^{2}\right\}\right] q d x \\
& =1+r \mathfrak{Y}+O\left(r^{2}\right), \\
\lim \log \mathfrak{M}_{r}= & \lim \frac{1}{r} \log \left\{1+r \mathfrak{Y}+O\left(r^{2}\right)\right\}=\mathfrak{Y} .
\end{aligned}
$$

There is some difficulty in extending this argument to the general case. The difficulty can however be avoided as follows ${ }^{\text {a }}$. By (6.7.6) and (6.7.7), we have

$$
\begin{align*}
\log (\mathscr{G}(f) & \leqq \log \mathfrak{M}_{r}(f)=\frac{1}{r} \log \mathfrak{A}\left(f^{r}\right)  \tag{6.8.2}\\
& \leqq \frac{1}{r}\left\{\mathfrak{A}\left(f^{r}\right)-1\right\}=\int \frac{f^{r}-1}{r} q d x .
\end{align*}
$$

When $r$ decreases to zero, $\left(t^{r}-1\right) / r$ decreases (by Theorem 36) and tends to the limit $\log t$. Hence ${ }^{\text {b }}$

$$
\begin{equation*}
\lim \frac{1}{r}\{\mathfrak{A}(f r)-1\}=\mathfrak{A}(\log f)=\log \mathscr{S}(f), \tag{6.8.3}
\end{equation*}
$$

the right-hand side being finite or $-\infty$. Combining (6.8.2) and (6.8.3) we see that

$$
\log \mathscr{S}(f) \leqq \underline{\lim } \log \mathfrak{M}_{r}(f) \leqq \varlimsup \overline{\lim } \log \mathfrak{M}_{r}(f) \leqq \log \mathfrak{G}(f),
$$

which proves the theorem.
6.9. Hölder's inequality for integrals. We consider next the
${ }^{\text {a }}$ F. Riesz (7). Other, less simple, proofs have been given by Besicovitch, Hardy, and Littlewood: see Hardy (7).
${ }^{\mathrm{b}}$ See § 6.3 (b) (ii).
theorems for integrals which correspond to Theorems 11-15. It is convenient to introduce another definition which enables us to shorten our statements of cases of equality. Two functions $f, g$ will be said to be effectively proportional if there are constants $A, B$, not both zero, such that $A f \equiv B g$. The idea has occurred already in Theorems 181 and 185. A nul function is effectively proportional to any function. We shall say that $f, g, h, \ldots$ are effectively proportional if every pair are so.
188. If $\alpha, \beta, \ldots, \lambda$ are positive and $\alpha+\beta+\ldots+\lambda=1$, then

$$
\begin{equation*}
\int f^{\alpha} g^{\beta} \ldots l^{\lambda} d x<\left(\int f d x\right)^{\alpha}\left(\int g d x\right)^{\beta} \ldots\left(\int l d x\right)^{\lambda} \tag{6.9.1}
\end{equation*}
$$

unless one of the functions is nul or all are effectively proportional.
Assuming no function nul, we have, by Theorem 9,

$$
\begin{aligned}
\frac{\int f^{\alpha} g^{\beta} \ldots l^{\lambda} d x}{\left(\int f d x\right)^{\alpha}\left(\int g d x\right)^{\beta} \ldots\left(\int l d x\right)^{\lambda}} & =\int\left(\frac{f}{\int f d x}\right)^{\alpha}\left(\frac{g}{\int g d x}\right)^{\beta} \ldots\left(\frac{l}{\int l d x}\right)^{\lambda} d x \\
& \leqq \int\left(\frac{\alpha f}{\int f d x}+\frac{\beta g}{\int g d x}+\ldots+\frac{\lambda l}{\int l d x}\right) d x=1
\end{aligned}
$$

with inequality unless

$$
\frac{f}{\int f d x} \equiv \frac{g}{\int g d x} \equiv \ldots \equiv \frac{l}{\int l d x} .
$$

As a corollary ${ }^{\text {a }}$ we have
189. If $k>1$ then

$$
\begin{equation*}
\int f g d x<\left(\int f^{k} d x\right)^{1 / k}\left(\int g^{k^{\prime}} d x\right)^{1 / k^{\prime}} \tag{6.9.2}
\end{equation*}
$$

unless $f^{k}$ and $g^{k^{\prime}}$ are effectively proportional.
If $0<k<1$ or $k<0$ then

$$
\begin{equation*}
\int f g d x>\left(\int f^{k} d x\right)^{1 / k}\left(\int g^{k^{\prime}} d x\right)^{1 / k^{\prime}} \tag{6.9.3}
\end{equation*}
$$

unless either (a) $f^{k}$ and $g^{k^{\prime}}$ are effectively proportional or (b) $f g$ is nul.

The second half of the theorem requires a little explanation. Suppose first that $0<k<1$ and that $\int g^{k^{\prime}} d x$ is finite, so that $g$ is almost always positive. If then we write $l=1 / k$, so that $l>1$, and

$$
f=(u v)^{l}, \quad g=v^{-l}
$$

so that

$$
\begin{gathered}
f g=u^{l}, \quad f^{k}=u v, \quad g^{k^{\prime}}=v^{l^{\prime}}, \\
\text { a Compare § 2.8. }
\end{gathered}
$$

then $u$ and $v$ are defined for almost all $x$, and
or

$$
\begin{gathered}
\int u v d x<\left(\int u^{l} d x\right)^{1 / l}\left(\int v^{l^{\prime}} d x\right)^{1 / l^{\prime}} \\
\int f^{k} d x<\left(\int f g d x\right)^{k}\left(\int g^{k^{\prime}} d x\right)^{1-k}
\end{gathered}
$$

unless $u^{l}, v^{l^{\prime}}$ are effectively proportional or, what is the same thing, unless $f^{k}, g^{k^{\prime}}$ are effectively proportional. Since $\int g^{k^{\prime}} d x$ is finite and not zero ${ }^{\text {a }}$, this is (6.9.3).

If $\int g^{k^{\prime}} d x=\infty$, then

$$
\left(\int g^{k^{\prime}} d x\right)^{1 / k^{\prime}}=0
$$

(since $k^{\prime}<0$ ). Hence the right-hand side of (6.9.3) is zero, and there is inequality unless $\int f g d x=0$, or $f g$ is nul.

When $k<0,0<k^{\prime}<1$, and the argument is substantially the same.

As we have explained in $\S \S 5.1$ and 6.5 , the theorem contains implicitly an assertion about convergence or finitude; if two of the integrals involved are finite, then so is the third. The integral which is finite if the other two are finite is $\int f g d x$ when $k>1, \int f^{k} d x$ if $0<k<1$, and $\int g^{k^{\prime}} d x$ if $k<0$.

The theorem corresponding to Theorem 161 is very important and, like Theorem 161, is not a direct corollary of preceding theorems.
190. ${ }^{\text {b }}$ If $k>1$ and fg belongs to $L$ for every $g$ which belongs to $L^{k^{\prime}}$, then $f$ belongs to $L^{k}$.

We consider first the case in which ( $a, b$ ) is finite (or $m E$ finite), and suppose that $\int f^{k} d x=\infty$. We can find a function $f^{*}$ which (1) has only an enumerable infinity of values $a_{i}$, and (2) satisfies $f^{*} \leqq f<f^{*}+\epsilon$. Since $f^{k}$ does not exceed a constant multiple of $f^{* k}+\left(f-f^{*}\right)^{k}$, by Theorem 13, we have $\int f^{* k} d x=\infty$. Hence, if $e_{i}$ is the set in which $f^{*}=a_{i}$,

$$
\Sigma a_{i}{ }^{k} e_{i}=\infty .
$$

It follows from Theorem 161, taking

$$
a_{i}{ }^{k} e_{i}=u_{i}{ }^{k}, \quad b_{i}{ }^{k^{\prime}} e_{i}=v_{i}{ }^{k^{\prime}}, \quad a_{i} b_{i} e_{i}=u_{i} v_{i},
$$

[^38]that there is a $b_{i}$ such that $\Sigma b_{i}{ }^{k^{\prime}} e_{i}$ is convergent and $\Sigma a_{i} b_{i} e_{\imath}=\infty$. We take $g(x)=b_{i}$ in $e_{i}$ (for all $i$ ). Then
$$
\int g^{k^{\prime}} d x=\Sigma b_{i}{ }^{k^{\prime}} e_{i}
$$
is convergent, but
$$
\int f^{*} g d x=\Sigma a_{i} b_{i} e_{i}=\infty
$$
and hence $\int f g d x=\infty$, contrary to the hypothesis.
If the integrals are over an infinite range, say ( $0, \infty$ ), we can write
$$
x=\frac{t}{1-t}
$$
when
$$
\int_{0}^{\infty} f g d x=\int_{0}^{1} F G d t, \quad \int_{0}^{\infty} f^{k} d x=\int_{0}^{1} F^{k} d t, \quad \int_{0}^{\infty} g^{k^{\prime}} d x=\int_{0}^{1} G^{k^{\prime}} d t
$$
where
$$
F(t)=(1-t)^{-2 / k} f\left(\frac{t}{1-t}\right), \quad G(t)=(1-t)^{-2 / k^{\prime}} g\left(\frac{t}{1-t}\right) .
$$

The theorem is thus reduced to the finite case.
191. If $k>1$, then a necessary and sufficient condition that $\int f^{k} d x \leqq F$ is that $\int f g d x \leqq F^{1 / k} G^{1 / k^{\prime}}$ for all $g$ such that $\int g^{k^{\prime}} d x \leqq G$.

The condition is necessary, by Theorem 189. If it is satisfied, then $\int f^{k} d x$ is finite, by Theorem 190. If $\int f^{k} d x>F$, we choose $g$ so that $g^{k^{\prime}}$ is effectively proportional to $f^{k}$, and then, by Theorem 189,

$$
\int f g d x=\left(\int f^{k} d x\right)^{1 / k}\left(\int g^{k^{\prime}} d x\right)^{1 / k^{\prime}}>F^{1 / k} G^{1 / k^{\prime}}
$$

The theorem may also be stated with ' $<$ ' for ' $\leqq$ ' in the first two inequalities: in order that $\int f^{k} d x<F$, it is necessary and sufficient that $\int f g d x<F^{1 / k} G^{1 / k^{\prime}}$ whenever $\int g^{k^{\prime}} d x \leqq G$.
We can prove Theorem 191 without appealing to the more difficult Theorem 190. If $\int f^{k} d x>F$ then $\int(f)_{n}{ }^{k} d x>F$ for sufficiently large $n$. Then, choosing $g$ effectively proportional to $\left(f_{n}\right)^{k-1}$, we have

$$
\int f g d x \geqq \int(f)_{n} g d x=\left(\int(f)_{n}^{k} d x\right)^{1 / k} G^{1 / k^{\prime}}>F^{1 / k} G^{1 / k^{\prime}},
$$

in contradiction to the hypothesis of the theorem.
Another proof of Theorem 190 (and of the associated Theorem 161) has been given by Banach (1, 85-86).

An example of the use of Theorem 191 appears in $\S 6.13$, in the proof of Theorem 202, and others in Ch. IX. ${ }^{\text {a }}$ In § 6.13 Theorem 202 is proved in
a See in particular $\S \S 9.3$ and 9.7 (2).
two different ways, of which one depends explicitly on Theorem 191 while the other does not use it, and the logical status of Theorem 191 in proofs of this character is explained in detail.
6.10. General properties of the means $\mathfrak{M}_{r}(f)$. We shallnow prove a number of theorems which include the analogues of those of §2.9. The properties to be investigated are a little more complex than they were there, and we shall require some additional conventions before we can state them comprehensively. We suppose first that $r>0$; the theorems which we prove in this case, with those which we have proved concerning $\mathfrak{G}(f)$, will give us the substance of what is required, and we shall be able to state the results for unrestricted $r$ more summarily, leaving most of the details of verification to the reader.
192. If $0<r<s$ and $\mathfrak{M}_{s}$ is finite, then

$$
\mathfrak{M}_{r}<\mathfrak{M}_{s},
$$

unless $f \equiv C$.
If $r=s \alpha$, so that $0<\alpha<1$, we have, by Theorem 188,

$$
\int q f^{r} d x=\int\left(q f^{s}\right)^{\alpha} q^{1-\alpha} d x<\left(\int q f^{s} d x\right)^{\alpha}\left(\int q d x\right)^{1-\alpha}=\left(\int q f^{s} d x\right)^{\alpha},
$$

unless $q f^{s} \equiv C q$. Since $q>0$, this is the result required.
193. If $\mathfrak{M}_{r}$ is finite for every $r$, then $\mathfrak{M}_{r} \rightarrow \operatorname{Max} f$ when $r \rightarrow+\infty$.
(i) Suppose $\mu=\operatorname{Max} f$ finite. Then (a) $\mathfrak{M}_{r} \leqq \mu$, and (b) $f>\mu-\epsilon$ in a set $e$ of positive measure $\zeta$, so that

$$
\int_{e} q d x=\zeta>0, \quad \mathfrak{M}_{r} \geqq(\mu-\epsilon) \zeta^{1 / r}, \quad \text { im } \mathfrak{M}_{r} \geqq \mu-\epsilon .
$$

(ii) Suppose $\operatorname{Max} f=\infty$. Then, for any $G>0, f>G$ in a set $e$ of positive measure, and, as above, $\lim \mathfrak{M}_{r} \geqq G$.

From (6.6.4), (6.6.5) and Theorem 193 it follows that

$$
\mathfrak{M}_{r} \rightarrow \operatorname{Min} f
$$

when $r \rightarrow-\infty$.
194. If $0<s<\infty$ and $\mathfrak{M}_{s}$ is finite, then $\mathfrak{M}_{r}$ is continuous for $0<r<$ s and continuous on the left for $r=s$. If $\mathfrak{M}_{s}=\infty$, but $\mathfrak{M}_{r}$ is finite for $0<r<s$, then $\mathfrak{M}_{r} \rightarrow \infty$ when $r \rightarrow s$.
(i) Suppose $\mathfrak{M}_{s}$ finite. Then

$$
q f^{r} \leqq q \operatorname{Max}\left(\mathbf{1}, f^{s}\right),
$$

a majorant of class $L$ independent of $r$; and the results follow from §6.3(b)(i). ${ }^{\text {a }}$
(ii) Suppose $\mathfrak{M}_{s}=\infty$. We can choose $n$ so that

$$
\int\left(q f^{s}\right)_{n} d x>G
$$

But $\left(q f^{r}\right)_{n}$ is a continuous function of $r$, and so ${ }^{b}$

$$
\int\left(q f^{r}\right)_{n} d x>\frac{1}{2} G
$$

for $r>s-\epsilon$. Hence $\int q f^{r} d x>\frac{1}{2} G$, which proves the theorem.
6.11. General properties of the means $\mathfrak{M}_{r}(f)$ (continued). In the preceding sectionswe have, in the main, confined our attention to means for which $r \geqq 0$, leaving it to the reader to deduce the corresponding results for means of negative order from the formulae (6.6.4) and (6.6.5). In this section we consider the means more comprehensively. We write, as is natural after Theorems 187 and 193,

$$
\begin{equation*}
\mathfrak{F}(f)=\mathfrak{M}_{0}(f), \operatorname{Max} f=\mathfrak{M}_{+\infty}(f), \operatorname{Min} f=\mathfrak{M}_{-\infty}(f) . \tag{6.11.1}
\end{equation*}
$$

$\mathfrak{M}_{0}(f)$ may be meaningless, but only if $\mathfrak{M}_{r}(f)=\infty$ for all $r>0$ and $\mathfrak{M}_{r}(f)=0$ for all $r<0$.

We begin by disposing of two exceptional cases.
(A) If $f \equiv C$ then $\mathfrak{M}_{r}=C$ for all $r$, and this is true even in the extreme cases $C=0$ and $C=\infty$.c
(B) We may have

$$
\mathfrak{M}_{r}=0(r<0), \quad \mathfrak{M}_{0} \text { meaningless, } \quad \mathfrak{M}_{r}=\infty(r>0) .
$$

These cases we dismiss. We then leave it to the reader to verify the truth of the assertions in (1) and (2) below, which cover all cases other than the exceptional cases (A) and (B).
(1) $\mathfrak{M}_{r}<\mathfrak{M}_{s}$ for $-\infty \leqq r<s \leqq \infty$, unless (a) $\mathfrak{M}_{r}=\mathfrak{M}_{s}=\infty$ (which can happen only if $r \geqq 0$ ), or (b) $\mathfrak{M}_{r}=\mathfrak{M}_{s}=0$ (which can happen only if $s \leqq 0$ ).

[^39](2) We denote by $\mathfrak{M}_{r-0}$ and $\mathfrak{M}_{r+0}$ the limits (which always exist) of $\mathfrak{M}_{t}$ when $t \rightarrow r$ from below and from above respectively.
If $r>0$ then $\mathfrak{M}_{r-0}=\mathfrak{M}_{r}$, and $\mathfrak{M}_{r+0}=\mathfrak{M}_{r}$ except when $\mathfrak{M}_{r}$ is positive and finite but $\mathfrak{M}_{t}=\infty$ for $t>r$, in which case
$$
\mathfrak{M}_{r+0}=\infty>\mathfrak{M}_{r}
$$

If $r<0$ then $\mathfrak{M}_{r+0}=\mathfrak{M}_{r}$, and $\mathfrak{M}_{r-0}=\mathfrak{M}_{r}$ except when $\mathfrak{M}_{r}$ is positive and finite but $\mathfrak{M}_{t}=0$ for $t<r$, in which case

$$
\mathfrak{M}_{r-0}=0<\mathfrak{M}_{r} .
$$

If $r=0$ there are exceptional cases corresponding to each of those indicated above. If $\mathfrak{M}_{0}$ is 0 or $\infty$ then either (a) $\mathfrak{M}_{-0}$ and $\mathfrak{M}_{+0}$ are each equal to $\mathfrak{M}_{0}$, or else (b)

$$
\mathfrak{M}_{-0}=\mathfrak{M}_{0}=0, \quad \mathfrak{M}_{+0}=\infty
$$

or $\quad \mathfrak{M}_{-0}=0, \quad \mathfrak{M}_{0}=\mathfrak{M}_{+0}=\infty$.
If $\mathfrak{M}_{0}$ is positive and finite, then each of $\mathfrak{M}_{-0}$ and $\mathfrak{M}_{+0}$, if also positive and finite, is equal to $\mathfrak{M}_{0}$; but $\mathfrak{M}_{-0}$ may also be 0 or $\mathfrak{M}_{+0}$ may be $\infty$.

Finally, all possibilities not explicitly excluded may actually occur ${ }^{\text {a }}$.

The results may be stated more symmetrically and concisely in terms of

$$
\mathfrak{E}_{r}=\log \mathfrak{M}_{r}:
$$

we agree that $\log \infty=+\infty$ and $\log 0=-\infty$. We put aside the cases corresponding to cases (A) and (B) above, viz.
(a) $f \equiv C$ (where $C$ may be 0 or $\infty$ ), when $\mathfrak{L}_{r}=\log C$ for all $r$;
(b) $\mathfrak{L}_{0}$ meaningless, when $\mathfrak{L}_{r}=+\infty$ for $r>0$ and $\mathfrak{L}_{r}=-\infty$ for $r<0$.
195. Apart from the cases just mentioned, the set of values of $r$ for which $\mathfrak{L}_{\boldsymbol{r}}=\log \mathfrak{M}$, is finite is either the nul set or a closed, half-closed, or open interval I or $(u, v)$, where $-\infty \leqq u \leqq v \leqq \infty$, which includes the point $r=0$ (so that $u \leqq 0 \leqq v$ ), but is otherwise arbitrary (so that, for example, $u$ may be $-\infty$ and $v$ be $+\infty$, or $u$ and $v$ may both be 0$). \mathfrak{L}_{\mathbf{r}}$ is $+\infty$ for values of $r$ to the right of $I$ and $-\infty$ for values to the left.

Inside $I, \mathfrak{L}_{r}$ is continuous and strictly increasing. If $r$ tends to an endpoint of I through values of $r$ interior to $I$, then $\mathfrak{L}_{r}$ tends to a limit (finite or infinite) equal to its value at the end-point in question.
6.12. Convexity of $\log \mathfrak{M}_{r}^{r}$. In this section (as in Theorem 17) we suppose $r>0$.

[^40]196. If $0<r<s<t$, and $\mathfrak{M}_{t}$ is finite, then
$$
\mathfrak{M}_{s}^{s}<\left(\mathfrak{M}_{r}^{r}\right)^{\frac{t-s}{t-r}}\left(\mathfrak{M}_{t}\right)^{\frac{s-r}{t-r}}
$$
unless $f \equiv 0$ in a part of $E$ and $f \equiv C$ in the complementary part.
The proof is based on Theorem 188, and is the analogue of that of Theorem 17. For equality,
$$
q f^{r} \equiv C q f^{t}
$$

As a corollary, we have
197. $\log \mathfrak{M}_{r}{ }^{r}(f)=r \log \mathfrak{M}_{r}(f)$ is a convex function of $r$.

Compare Theorem 87. The reader will find it instructive to deduce the continuity of $\mathfrak{M}_{r}$ (Theorem 194) from Theorem 197.
6.13. Minkowski's inequality for integrals. The inequalities of the Minkowski type are derived in substantially the same way as in $\S 2.11$. The ordinary form of Minkowski's inequality for integrals is
198. If $k>1$ then
(6.13.1) $\left\{\int(f+g+\ldots+l)^{k} d x\right\}^{1 / k}<\left(\int f^{k} d x\right)^{1 / k}+\ldots+\left(\int l^{k} d x\right)^{1 / k}$, and if $0<k<1$ then
(6.13.2) $\quad\left\{\int(f+g+\ldots+l)^{k} d x\right\}^{1 / k}>\left(\int f^{k} d x\right)^{1 / k}+\ldots+\left(\int l^{k} d x\right)^{1 / k}$, unless $f, g, \ldots, l$ are effectively proportional.

The inequality (6.13.2) is still true generally when $k<0$, but there is a second case of exception, when both sides of the inequality vanish.

We deduce this from Theorem 189 much as we deduced Theorem 24 from Theorem 13. Since the cases of equality are a little puzzling, we write out the proof of (6.13.2) in detail.

If $S=f+g+\ldots+l$ then
(6.13.3) $\quad \int S^{k} d x=\int f S^{k-1} d x+\int g S^{k-1} d x+\ldots+\int l S^{k-1} d x$.

Suppose first that $0<k<1$. By Theorem 19,

$$
S^{k} \leqq f^{k}+g^{k}+\ldots+l^{k}
$$

Hence, if $\int f^{k} d x, \ldots$ are finite, $\int S^{k} d x$ is finite. Also $\int S^{k} d x>0$
unless $S \equiv 0$ and so $f, g, \ldots$ are all nul. We may therefore suppose $\int S^{k} d x$ positive and finite.

By Theorem 189,

$$
\int f S^{k-1} d x>\left(\int f^{k} d x\right)^{1 / k}\left(\int S^{k} d x\right)^{1 / k^{\prime}},
$$

unless (a) $f^{k}$ and $S^{k}$ are effectively proportional or (b) $f S^{k-1} \equiv 0$. Since $k-1<0$, and $S$ is finite almost everywhere, the second alternative can occur only if $f$ is nul, and so reduces to a case of the first. Hence (6.13.3) gives
(6.13.4) $\int S^{k} d x>\left\{\left(\int f^{k} d x\right)^{1 / k}+\ldots+\left(\int l^{k} d x\right)^{1 / k}\right\}\left(\int S^{k} d x\right)^{1 / k^{\prime}}$,
unless $f, g, \ldots, l$ are effectively proportional; and the conclusion follows.

The argument goes similarly when $k<0$, provided $\int S^{k} d x$ is positive and finite. If $\int S^{k} d x=0$ then, since $k<0, S$ is infinite almost everywhere, which is impossible since every $f$ is finite almost everywhere. If $\int S^{k} d x$ is infinite then (again since $k<0$ ) $\int f^{k} d x, \ldots$ are all infinite, and both sides of (6.13.2) are zero. This is the second exceptional case mentioned in the enunciation, and occurs, for example, when

$$
f=g=\ldots=l=0
$$

in a set $E$ of positive measure.
We have excluded the cases $k=1$ and $k=0$ from the statement of Theorem 198. The first is trivial and the second is included in Theorem 186. We leave it to the reader to state Theorem 198 in a form corresponding to that of Theorem 24.

Corresponding to Theorem 27, we have
199. If $k>1$ then
(6.13.5) $\int(f+g+\ldots+l)^{k} d x>\int f^{k} d x+\ldots+\int l^{k} d x$, and if $0<k<1$ then
(6.13.6) $\int(f+g+\ldots+l)^{k} d x<\int f^{k} d x+\ldots+\int l^{k} d x$, unless, for almost all $x$, all but one of $f, g, \ldots, l$ are zero. If all of $f, g, \ldots, l$ are almost always positive, then (6.13.6) is true also for $k<0$.

Theorem 198, with $k>1$, is a special case of the first of the following three more general theorems, in which the series are
finite or infinite and the ranges of integration arbitrary. We confine ourselves to the case $k>1$; in general, the sign of inequality is reversed when $k<1$.
200. If $k>1$ then

$$
\begin{equation*}
\left[\int\left\{\Sigma f_{m}(x)\right\}^{k} d x\right]^{1 / k}<\Sigma\left\{\int f_{m}^{k}(x) d x\right\}^{1 / k}, \tag{6.13.7}
\end{equation*}
$$

unless

$$
f_{m}(x) \equiv C_{m} \phi(x)
$$

201. If $k>1$ then
(6.13.8)
unless

$$
\begin{gathered}
{\left[\Sigma\left\{\int f_{n}(x) d x\right\}^{k}\right]^{1 / k}<\int\left\{\Sigma f_{n}^{k}(x)\right\}^{1 / k} d x} \\
f_{n}(x) \equiv C_{n} \phi(x) .
\end{gathered}
$$

202. If $k>1$ then
(6.13.9) $\left[\int\left\{\int f(x, y) d y\right\}^{k} d x\right]^{1 / k}<\int\left\{\int f^{k}(x, y) d x\right\}^{1 / k} d y$, unless $f(x, y) \equiv \phi(x) \psi(y)$.
In each theorem there is equality in the exceptional case.
Consider for example Theorem 202 (the least elementary of the theorems). We begin by proving the theorem with ' $\leqq$ '. We give two proofs, in the first of which we appeal to Theorem 191. In each proof the chain of equalities and inequalities which arises is to be interpreted in the sense 'if the right-hand side of any equality or inequality is finite, then so is the left-hand side, and the two are related as stated'. The inversions of the order of integration are justified by Fubini's Theorem.

We write

$$
J=J(x)=\int f(x, y) d y
$$

(i) In order that

$$
\begin{equation*}
\int J^{k} d x \leqq M^{k} \tag{6.13.10}
\end{equation*}
$$

it is, by Theorem 191, necessary and sufficient that

$$
\int J g d x \leqq M
$$

for all $g$ for which

$$
\begin{equation*}
\int g^{k^{\prime}} d x \leqq 1 \tag{6.13.12}
\end{equation*}
$$

Now
(6.13.13) $\quad \int J g d x=\int g(x) d x \int f(x, y) d y$

$$
=\int d y\left(\int g(x) f(x, y) d x\right) \leqq \int d y\left(\int f^{k}(x, y) d x\right)^{1 / k}
$$

by Theorem 189 and (6.13.12). Hence we may take

$$
M=\int d y\left(\int f^{k} d x\right)^{1 / k}
$$

in (6.13.10), which proves the theorem (with ' $\leqq$ ').
(ii) If $\int J^{k} d x=0$, then $J=0$ for almost all $x$, and so (for almost all $x) f=0$ for almost all $y$. Hence, after $\S 6.3(d), f(x, y) \equiv 0$.

We may therefore suppose that $\int J^{k} d x>0$. Let us assume for a moment that $\int J^{k} d x$ is finite. Then

$$
\begin{aligned}
\int J^{k} d x=\int J^{k-1} d x \int f d y & =\int d y \int J^{k-1} f d x \leqq \int d y\left\{\left(\int f^{k} d x\right)^{1 / k}\left(\int J^{k} d x\right)^{1 / k^{\prime}}\right\} \\
& =\left(\int J^{k} d x\right)^{1 / k^{\prime}} \int\left(\int f^{k} d x\right)^{1 / k} d y
\end{aligned}
$$

and so

$$
\begin{equation*}
\left(\int J^{k} d x\right)^{1 / k} \leqq \int\left(\int f^{k} d x\right)^{1 / k} d y \tag{6.13.14}
\end{equation*}
$$

which is (6.13.9), with ' $\leqq$ ' for ' $<$ '.
In this proof we have assumed $\int J^{k} d x$ finite, an assumption which was not required in proof (i). In order to get rid of the assumption, we must approximate to $f$ by some function for which the assumption is certainly justified. Suppose for example that the integrations are over finite intervals or sets of finite measure, that $(f)_{n}$ is defined as in $\S 6.1$, and that

$$
J_{n}=\int(f)_{n} d y
$$

Then $\int J_{n}{ }^{k} d x$ is certainly finite, and so

$$
\left(\int J_{n}^{k} d x\right)^{1 / k} \leqq \int\left\{\int(f)_{n}^{k} d x\right\}^{1 / k} d y \leqq \int\left(\int f^{k} d x\right)^{1 / k} d y
$$

From this (6.13.9) follows, with ' $\leqq$ ', by making $n \rightarrow \infty$.
The arguments under (i) and (ii) are essentially of the same character, the part of the arbitrary $g$ in (i) being played in (ii) by the definite function

$$
g=\frac{J^{k-1}}{\left(\int J^{k} d x\right)^{1 / k^{\prime}}}
$$

which satisfies (6.13.12) if $\int J^{k} d x$ is finite. By using this particular $g$, we avoid an appeal to a rather sophisticated general theorem, but at the cost of some additional complications. A similar alternative presents itself whenever we make use of Theorem 191.

It remains to discuss the possibility of equality in (6.13.9). There will be inequality if ${ }^{\text {a }}$

$$
\int J g d x<M
$$

for all $g$ subject to (6.13.12). There is inequality in

$$
\int d y\left(\int g f d x\right) \leqq \int d y\left\{\left(\int f^{k} d x\right)^{1 / k}\left(\int g^{k^{\prime}} d x\right)^{1 / k^{\prime}}\right\}
$$

a See the last remark of $\S 6.9$.
unless, for almost all $y, f^{k}$ and $g^{k^{\prime}}$ are effectively proportional, i.e. unless (for almost all $y$ )

$$
\begin{equation*}
\rho(y) f^{k}(x, y)=\sigma(y) g^{k^{\prime}}(x) \tag{6.13.15}
\end{equation*}
$$

where $\rho^{2}+\sigma^{2}>0$, for almost all $x$. If $\rho(y)$ were zero, for a $y$ for which (6.13.15) holds, $g(x)$ would be nul, which is false. Hence, in (6.13.15), $\rho(y)>0$, and so

$$
f(x, y)=\phi(x) \psi(y)
$$

where $\phi=g^{k^{\prime} / k}, \psi=(\sigma / \rho)^{1 / k}$. This equation holds, for almost all $y$, for almost all $x$, and therefore, by $\S 6.3(d)$, for almost all $x, y$.

The proofs of Theorems 200 and 201 follow similar lines. Thus, in proving Theorem 201, we write

$$
J_{n}=\int f_{n} d x
$$

and argue as follows. In order that $\Sigma J_{n}{ }^{k}<M^{k}$, it is necessary and sufficient, by Theorem 15, a that $\Sigma b_{n} J_{n}<M$ whenever $\Sigma b_{n}{ }^{k \prime} \leqq 1$. Also

$$
\Sigma b_{n} J_{n}=\Sigma b_{n} \int f_{n} d x=\int\left(\Sigma b_{n} f_{n}\right) d x \leqq \int d x\left(\Sigma f_{n}^{k}\right)^{1 / k}\left(\Sigma b_{n}^{k^{\prime}}\right)^{1 / k^{\prime}} \leqq \int\left(\Sigma f_{n}^{k}\right)^{1 / k} d x
$$

and so on. The summation under the integral sign is justified by (ii) of §6.3(b).

The analogue of Theorem 26 is
203. If $0<r<s$ then

$$
\mathfrak{M}_{s}{ }^{(y)} \mathfrak{M}_{r}{ }^{(x)} f(x, y)<\mathfrak{M}_{r}{ }^{(x)} \mathfrak{M}_{s}{ }^{(y)} f(x, y),
$$

unless $f(x, y) \equiv \phi(x) \psi(y)$.
For an explicit proof see Jessen (1).
6.14. Mean values depending on an arbitrary function. There is a theory of integral mean values involving an arbitrary function similar to that developed in Ch. III. We do not set it out in detail here because it would be so largely a repetition, in a slightly different form, of what we have said already. We confine ourselves to proving the analogue of Theorem 95. ${ }^{\text {b }}$
204. Suppose that $\alpha \leqq f(x) \leqq \beta$, where $\alpha$ and $\beta$ may be finite or infinite, and that $f(x)$ is almost always different from $\alpha$ and $\beta$; that the range of integration and the weight function $p(x)$ satisfy the

[^41]conditions of § 6.6; and that $\phi^{\prime \prime}(t)$ is positive and finite for $\alpha<t<\beta$. Then
\[

$$
\begin{equation*}
\phi\left(\frac{\int f p d x}{\int p d x}\right) \leqq \frac{\int \phi(f) p d x}{\int p d x}, \tag{6.14.1}
\end{equation*}
$$

\]

whenever the right-hand side exists and is finite; and there is equality only when $f \equiv C$.

It is possible that $\int f p d x=\infty$ or $\int f p d x=-\infty ;(6.14 .1)$ is then still true if interpreted in the obvious manner. It is not possible (when the right-hand side is finite) that $\int f p d x$ should not exist, i.e. $\int f^{+} p d x=\infty$ and $\int f^{-} p d x=-\infty$. For in this case $\alpha=-\infty$, $\beta=\infty$, and $\phi(f)$, being convex and not constant, must tend to infinity, with rapidity at least that of a multiple of $|f|$, either for large positive or for large negative values of $f,{ }^{\text {a }}$ so that $\int \phi(f) p d x$ cannot exist and be finite.

We take $p=q, \int q d x=1$, and suppose first that $\mathfrak{M}=\int f q d x$ is finite. If $f$ is not effectively constant, $\alpha<\mathfrak{M}<\beta$. Also $f$ is finite, and $\alpha<f<\beta$, for almost all $x$; so that, for almost all $x$,

$$
\phi(f)=\phi(\mathfrak{M})+(f-\mathfrak{M}) \phi^{\prime}(\mathfrak{M})+\frac{1}{2}(f-\mathfrak{M})^{2} \phi^{\prime \prime}(\mu),
$$

where $\mu$ lies between $f$ and $\mathfrak{M}$, so that $\alpha<\mu<\beta$. Hence

$$
\int \phi(f) q d x \geqq \phi(\mathfrak{M}),
$$

which is (6.14.1). There is equality only if $(f-\mathfrak{M} \mathcal{C})^{2} \phi^{\prime \prime}(\mu) \equiv 0$; but $\alpha<\mu<\beta$, and so $\phi^{\prime \prime}(\mu)>0$, for almost all $x$, so that then $f \equiv \mathbb{M}$.

Next suppose (say) $\int f q d x=\infty$, so that $\beta=\infty$. Then

$$
\phi\left\{\int(f)_{n} q d x\right\} \leqq \int \phi\left\{(f)_{n}\right\} q d x
$$

by what has been proved already. Since $\phi(f)$ is continuous and monotonic for large $f$, the integral on the right tends to $\int \phi(f) q d x$, while that on the left tends to $\phi(\infty)$. Hence $\phi(\infty)$ is finite, in which case $\phi$ is decreasing and

$$
\phi(\infty)<\phi(f)
$$

It follows that

$$
\phi(\infty)=\phi(\infty) \int q d x \leqq \int \phi(f) q d x
$$

with equality only if $\phi(f) \equiv \phi(\infty)$, a possibility which we excluded. The case in which $\int f q d x=-\infty$ may be discussed similarly.
a See Theorem 126.

It is possible that the left-hand side of (6.14.1) should be $-\infty$. The reader will find it instructive to verify that the various cases which we have contemplated can all occur.

If we take $\phi(t)=-\log t$, we obtain

$$
\exp \left(\int q \log f d x\right) \leqq \int q f d x
$$

i.e. $\left(\$(f) \leqq \mathfrak{A}(f)\right.$ (Theorem 184). If we take $\phi=t^{r}$, we are led again to Hölder's inequality, and other examples may be constructed analogous to those of $\S 3.11$. If we take $\phi(t)=t \log t$, we find
205. $\frac{\int p f d x}{\int p d x}<\exp \left(\frac{\int p f \log f d x}{\int p f d x}\right)$,
unless $f \equiv C$.
We can extend the result of Theorem 204 (except for the specification of the cases of equality) to any convex and continuous $\phi$.
206. The inequality (6.14.1) is true whenever $\phi(t)$ is convex and continuous in $\alpha<t<\beta$.

After § 3.19, we have

$$
\phi(f) \geqq \phi(\mathfrak{M})+\lambda(f-\mathfrak{M}),
$$

where $\lambda$ is any number between the left and right hand derivatives of $\phi(t)$ for $t=\mathfrak{M}$. Hence

$$
\int \phi(f) q d x \geqq \phi(\mathfrak{M}),
$$

which is (6.14.1).

## STIELTJES INTEGRALS

6.15. The definition of the Stieltjes integral. We have so far considered series and integrals separately, and all the fundamental theorems have appeared in dual form; thus Hölder's inequality is contained in Theorems 13 and 189. It is natural to look for an extension of these theorems which combines them into one, and we can find such an extension by using Stieltjes integrals.

Suppose that $\phi(x)$ increases (in the wide sense) in $a \leqq x \leqq b$, and
that $\phi(a)=\alpha, \phi(b)=\beta$. We suppose $\alpha$ and $\beta$ (but not necessarily $a$ and $b)$ finite ${ }^{\text {a }}$. The curve

$$
y=y(x)=\phi(x)
$$

is a rising curve which may have an enumerable set of ordinary discontinuities or of stretches of invariability. The inverse function

$$
x=x(y)=x(\phi)
$$

is defined uniquely except $(a)$ in intervals $\left(y_{1}, y_{2}\right)$ of $y$ corresponding to discontinuities $x=\xi$ of $\phi$ and (b) for values of $y$ which correspond to stretches of invariability of $\phi$. If we agree that $\left(y_{1}, y_{2}\right)$ is a stretch of invariability of $x(y)$, in which it has the value $\xi$, then $x(y)$ is defined except for the values $(b)$, and is an increasing function of $y$ for the values of $y$ for which it is defined. Finally we complete the definition of $x(y)$, as an increasing function of $y$, by assigning to it, for a value (b) of $y$, any one of the values of $x$ in the stretch of invariability. These values of $y$ are enumerable, and our choice of $x(y)$ for any of them has no effect on the definitions which follow.

We now define the Stieltjes, or Lebesgue-Stieltjes, integral

$$
\int_{x=a}^{x=b} f(x) d \phi(x)=\int_{a}^{b} f(x) d \phi
$$

of $f(x)$ with respect to $\phi(x)$, by

$$
\begin{equation*}
\int_{a}^{b} f(x) d \phi=\int_{a}^{\beta} f\{x(\phi)\} d \phi \tag{6.15.1}
\end{equation*}
$$

whenever the integral on the right-hand side exists as a Lebesgue integral ${ }^{\text {b }}$.

The definition (6.15.1), due to Radon (1), reduces the theory of Stieltjes integrals to that of Lebesgue integrals, and we may therefore expect that no new difficulties will arise. For full discussions of this and older definitions of the Stieltjes integral, we may refer to Hobson (1), Lebesgue (1), Pollard (1), Young (7).

[^42]We shall not require this more general definition here.

We can define

$$
\int_{E} f(x) d \phi
$$

where $\phi$ is an increasing function, and $E$ a set of values of $x$, similarly, that is to say by the equation

$$
\int_{E} f(x) d \phi=\int_{\varepsilon} f\{x(\phi)\} d \phi,
$$

where $\varepsilon$ is the set of values of $\phi$ corresponding to $E$. We must assume $\varepsilon$ measurable. The integral
is the variation of $\phi$ in $E$.

$$
\int_{E} d \phi
$$

6.16. Special cases of the Stieltjes integral. The simplest cases are the following:
(a) $\phi=x$. In this case the Stieltjes integral reduces to the ordinary Lebesgue integral.
(b) $\phi$ is an integral. In this case

$$
\int_{a}^{b} f(x) d \phi=\int_{a}^{b} f(x) \phi^{\prime}(x) d x .
$$

(c) $\phi$ is a finite increasing step-function.

Suppose that $a=a_{1}<a_{2}<\ldots<a_{n}=b$, that $\phi(x)=\alpha_{k}$, where $\alpha_{k}<\alpha_{k+1}$, in $a_{k}<x<a_{k+1}$, and that $\phi\left(a_{k}\right)$, when $1<k<n$, has any value consistent with the fact that $\phi$ increases. Then $x(y)$ is a step-function with values $a_{1}, a_{2}, \ldots, a_{n}$, and

$$
\begin{align*}
& \text { 1) } \int_{a}^{b} f d \phi=\int_{a}^{\beta} f\{x(\phi)\} d \phi  \tag{6.16.1}\\
& =\left(\alpha_{1}-\alpha\right) f\left(a_{1}\right)+\left(\alpha_{2}-\alpha_{1}\right) f\left(a_{2}\right)+\ldots+\left(\alpha_{n-1}-\alpha_{n-2}\right) f\left(a_{n-1}\right) \\
& =\Sigma \rho_{k} f\left(a_{k}\right), \\
& +\left(\beta-\alpha_{n-1}\right) f\left(a_{n}\right)
\end{align*}
$$

where $\rho_{k}$ is the saltus of $\phi$ at $x=a_{k}$. It is plain that any finite sum can be expressed as a Stieltjes integral; thus

$$
\sum_{1}^{n} u_{k}=\int f(x) d \phi
$$

where $\phi$ is a step-function with unit jumps at $a_{1}, a_{2}, \ldots, a_{n}$, and $u_{k}=f\left(a_{k}\right)$.
(d) These considerations extend at once to step-functions with infinitely many discontinuities, when the Stieltjes integral is $\Sigma \rho_{k} f\left(a_{k}\right)$, summed over all the discontinuities. Any convergent infinite series may be expressed in this way as a Stieltjes integral.
6.17. Extensions of earlier theorems. It will now be clear that all our fundamental theorems may be extended at once to Stieltjes integrals, and that the theorems thus obtained include those for Lebesgue integrals and also those for sums. We state the most representative of these theorems in the next section. Two preliminary remarks will be useful.
(1) When the Stieltjes integral is written as a Lebesgue integral, the variable of integration is $\phi$. Our conditions for equality were always of the type $f \equiv g, f=g$ except in a set of measure zero. The exceptional set in our new theorems will be of measure zero in $\phi$, and when we state this concept again in terms of $x$ it becomes 'a set of values of $x$ in which the variation of $\phi$ is zero', i.e. a set $E$ such that the corresponding values of $\phi$ form a nul set. Our conditions for equality must therefore all be interpreted in this sense. Thus ' $f$ is effectively proportional to $g$ ' means that

$$
A f=B g
$$

where $A$ and $B$ are constants, not both zero, except at the points of a set over which the variation of $\phi$ is zero. It will be observed that such an exceptional set cannot include any point at which $\phi(x)$ is discontinuous.

A similar point occurs in the definition of $\operatorname{Max} f$ and $\operatorname{Min} f$. Thus Max $f$ is the greatest number $\xi$ such that, for every positive $\epsilon$, $f>\xi-\epsilon$ in a set in which the variation of $\phi$ is positive.
(2) Many inequalities ' $X<Y$ ' are true for Lebesgue integrals when their analogues for Stieltjes integrals are true only with ' $\leqq$ '. Suppose, for example, that the integrations are over ( $0, \infty$ ) and that $\int f d x=1$. Then, by Theorem 181,

$$
\begin{equation*}
\left(\int x f d x\right)^{2}<\int f d x \int x^{2} f d x=\int x^{2} f d x \tag{6.17.1}
\end{equation*}
$$

unless $x^{2} f \equiv C f$ or $x^{2} \equiv C$, which is untrue, so that (6.17.1) is true
in any case. In the corresponding theorem for Stieltjes integrals we have $\int d \phi=1$ and

$$
\begin{equation*}
\left(\int x d \phi\right)^{2} \leqq \int d \phi \int x^{2} d \phi=\int x^{2} d \phi \tag{6.17.2}
\end{equation*}
$$

There is equality in (6.17.2) if $x^{2} \equiv C$, i.e. if $x$ is constant except in a set over which the variation of $\phi$ is zero, or, what is the same thing, if $\phi$ varies at one point only. Thus if $\phi=0$ for $0 \leqq x<1$, and $\phi(x)=1$ for $x \geqq 1$, then

$$
\left(\int x d \phi\right)^{2}=1=\int x^{2} d \phi
$$

6.18. The means $\mathfrak{M}_{r}(f ; \phi)$. We write

$$
\begin{aligned}
\mathfrak{M}_{r}(f)=\mathfrak{M}_{r}(f ; \phi)= & \left(\frac{1}{\beta-\alpha} \int_{a}^{b} f^{r} d \phi\right)^{1 / r}=\left(\frac{\int f^{r} d \phi}{\int d \phi}\right)^{1 / r}(r \neq 0), \\
& \mathfrak{U}(f ; \phi)=\mathfrak{M}_{\mathbf{1}}(f ; \phi), \\
\mathfrak{F}(f ; \phi)= & \exp \left(\frac{\int \log f d \phi}{\int d \phi}\right)=\mathfrak{M}_{0}(f ; \phi) .
\end{aligned}
$$

These definitions presuppose that the integrals involved are finite. If $\int f^{r} d \phi=\infty$, we agree (following the conventions of §6.6) that $\mathfrak{M l}_{r}=\infty$ when $r>0$ and $\mathfrak{M}_{r}=0$ when $r<0$. The points discussed in $\S \S 5.2$ and 6.7 naturally recur here in connection with the definition of 85 .

The theorems corresponding to Theorems 183, 184, 187, 189, 192, 193, 197, and 198 are as follows: we suppose for simplicity of statement that $r>0$.
207. $\operatorname{Min} f<\mathfrak{M}_{r}(f)<\operatorname{Max} f$ unless $f \equiv C$.
208. $\mathfrak{G}(f)<\mathfrak{M}_{r}(f)$, and in particular $\mathfrak{H}(f)<\mathfrak{H}(f)$, unless $f \equiv C$.
209. If $\mathfrak{M}_{r}(f)$ is finite for some $r$, then $\mathfrak{M}_{r}(f) \rightarrow \mathscr{G}(f)$ when $r \rightarrow+0$.
210. If $k>1$, then

$$
\int u v d \phi<\left(\int u^{k} d \phi\right)^{1 / k}\left(\int v^{k^{\prime}} d \phi\right)^{1 / k^{\prime}}
$$

unless $u^{k}$ and $v^{k^{\prime}}$ are effectively proportional. The inequality is reversed when $0<k<1$ or $k<0$, except when $u^{k}$ and $v^{k^{\prime}}$ are effectively proportional, or the left-hand side is zero (in which case the right-hand side is also zero).

This is Hölder's inequality; there are naturally corresponding generalisations of Theorems 11 (or 10) and 188.
211. If $r<s$, then $\mathfrak{M}_{r}(f)<\mathbb{M}_{s}(f)$, unless $f \equiv C$.
212. If $\mathfrak{M}_{r}(f)$ is finite for every positive $r$, then $\mathfrak{M}_{r}(f) \rightarrow \operatorname{Max} f$ when $r \rightarrow+\infty$.
213. $\log \mathfrak{M}_{r}^{r}(f)$ is a convex function of $r$.
214. If $k>1$, then

$$
\left\{\int(u+v)^{k} d \phi\right\}^{1 / k}<\left(\int u^{k} d \phi\right)^{1 / k}+\left(\int v^{k} d \phi\right)^{1 / k},
$$

unless $u$ and $v$ are effectively proportional. The inequality is in general reversed if $0<k<1$ or $k<0$. ${ }^{\text {a }}$

## AXIOMATIC TREATMENT OF MEAN VALUES

6.19. Distribution functions. In Ch. III we defined the mean value

$$
\mathfrak{M}_{\phi}=\mathfrak{M}_{\phi}(a, q)=\phi^{-1}\{\Sigma q \phi(a)\}
$$

directly, and developed its characteristic properties from the definition. Here we reverse the process and give the 'axiomatic' treatment promised on p.66. It is convenient to use the notation of Stieltjes integration, and it is for this reason that we have reserved the discussion until now; but the Stieltjes integrals which we use are actually all finite sums.

In what follows we consider a special class of step-functions, defined for all real $x$, which we call finite distribution functions. We call $F(x)$ a finite distribution function if
(i) it is constant in stretches and has only a finite number of discontinuities,
(ii) it increases (in the wide sense) from 0 to 1 , so that

$$
F(-\infty)=0, \quad F(\infty)=1,
$$

(iii) $F(x)=\frac{1}{2}\{\boldsymbol{F}(x-0)+F(x+0)\}$ for all $x$.

The distribution function which has jumps $q$ at the points $a$ provides a representation of both the values $a$ and the weights $q$ involved in $\mathfrak{M r}_{\phi}(a)$. The simplest such function is

$$
E(x)=\frac{1}{2}(1+\operatorname{sgn} x),
$$

[^43]which has the single jump 1 at $x=0$. If we write
$$
E_{\xi}(x)=E(x-\xi)
$$
then
\[

$$
\begin{equation*}
F(x)=\Sigma q E_{a}(x) \tag{6.19.1}
\end{equation*}
$$

\]

where

$$
a=a_{\nu}, \quad q=q_{\nu}(\nu=1,2, \ldots, n), \quad \Sigma q_{\nu}=1, \quad a_{1}<a_{2}<\ldots<a_{n}
$$

is the general finite distribution function with jumps $q$ at the points a. Also

$$
\begin{equation*}
\int_{-\infty}^{\infty} \phi(x) d F(x)=\Sigma q \phi(a) \tag{6.19.2}
\end{equation*}
$$

and the mean value (3.1.3) may be written as

$$
\begin{equation*}
\mathfrak{M}_{\phi}[F]=\phi^{-1}\left(\int_{-\infty}^{\infty} \phi(x) d F(x)\right) \tag{6.19.3}
\end{equation*}
$$

Any finite distribution function is 0 for $x<A$ and 1 for $x>B$, $A$ and $B$ being finite numbers depending on $F$. In what follows we confine our attention to a sub-class of these functions, viz. those which satisfy

$$
\begin{equation*}
F(x)=0 \quad(x<A), \quad F(x)=1 \quad(x>B) \tag{6.19.4}
\end{equation*}
$$

for a fixed $A$ and $B$. In these circumstances we say that $F$ belongs to $\mathfrak{D}(A, B)$.

If $\phi(x)$ is continuous and strictly monotonic in the closed interval $(A, B)$, then $\mathfrak{M}_{\phi}[F]$ is defined, by (6.19.3), for all $F$ of $\mathscr{D}(A, B)$. The values of $\phi(x)$ outside $(A, B)$ are not really involved in (6.19.3), and we may choose them as we please; it is natural to choose them so that $\phi(x)$ is continuous and strictly monotonic for $-\infty \leqq x \leqq \infty$.
6.20. Characterisation of mean values. Our object is to prove the following theorem.
215. Suppose that there is a unique real number $\mathfrak{M}[F]$, corresponding to each $F$ of $\mathfrak{D}(A, B)$, with the following properties:

$$
\begin{equation*}
\mathfrak{M}\left[E_{\xi}(x)\right]=\xi \quad(A \leqq \xi \leqq B) ; \tag{1}
\end{equation*}
$$

[2] if $F_{1}$ and $F_{2}$ belong to $\mathfrak{D}(A, B), F_{1} \geqq F_{2}$ for all $x$, and $F_{1}>F_{2}$ for some $x$, then

$$
\mathfrak{M d}\left[F_{1}\right]<\mathfrak{M}\left[F_{2}\right] ;
$$

[3] if $F, F^{*}, G$ belong to $\mathfrak{D}(A, B)$, and

$$
\mathfrak{M}[F]=\mathfrak{M}\left[F^{*}\right],
$$

then

$$
\mathfrak{M}[t F+(1-t) G]=\mathfrak{M}\left[t F^{*}+(1-t) G\right]
$$

for $0<t<1$.
Then there is a function $\phi(x)$, continuous and strictly increasing in the closed interval $(A, B)$, for which

$$
\begin{equation*}
\mathfrak{M}[F]=\mathfrak{M}_{\phi}[F]=\phi^{-1}\left(\int_{-\infty}^{\infty} \phi(x) d F(x)\right) . \tag{6.20.1}
\end{equation*}
$$

Conversely, if $\mathfrak{M}[F]$ is defined by (6.20.1), for a $\phi(x)$ with the properties stated, then it satisfies [1], [2], and [3], so that these conditions are necessary and sufficient for the representation of $\mathfrak{M}[F]$ in the form (6.20.1). ${ }^{\text {a }}$

We begin by proving the converse half of the theorem. If $\mathfrak{M}[F]$ is defined by (6.20.1), then it is obvious that it possesses property [1], and all but obvious that it possesses property [3], since

$$
\begin{aligned}
\phi(\mathfrak{M}[t F+(1-t) G]) & =t \int \phi d F+(1-t) \int \phi d G \\
& =t \int \phi d F^{*}+(1-t) \int \phi d G=\phi\left(\mathfrak{M}\left[t F^{*}+(1-t) G\right]\right) .
\end{aligned}
$$

It remains to prove [2].
Suppose that $F_{1}$ and $F_{2}$ satisfy the conditions stated. Then there is a positive number $\mu$ and an interval $(\alpha, \beta)$ such that

$$
F_{1}(x)>F_{2}(x)+\mu>F_{2}(x)
$$

in $(\alpha, \beta) \cdot{ }^{\text {b }}$ Hence

$$
\begin{aligned}
\phi\left(\mathfrak{M}\left[F_{2}\right]\right)-\phi\left(\mathfrak{M}\left[F_{1}\right]\right) & =\int_{-\infty}^{\infty} \phi d F_{2}-\int_{-\infty}^{\infty} \phi d F_{1} \\
& =\int_{-\infty}^{\infty}\left(F_{1}-F_{2}\right) d \phi^{\mathrm{c}} \\
& \geqq \int_{\alpha}^{\beta}\left(F_{1}-F_{2}\right) d \phi \geqq \mu\{\phi(\beta)-\phi(\alpha)\}>0 .
\end{aligned}
$$

a See Nagumo (1), Kolmogoroff (1), de Finetti (1). We follow the lines of de Finetti's proof.
b There is an $x_{0}$ for which $F_{1}\left(x_{0}\right)>F_{2}\left(x_{0}\right)$ or

$$
\frac{1}{2}\left\{F_{1}\left(x_{0}-0\right)+F_{1}\left(x_{0}+0\right)\right\}>\frac{1}{2}\left\{F_{2}\left(x_{0}-0\right)+F_{2}\left(x_{0}+0\right)\right\} .
$$

Hence either $F_{1}\left(x_{0}-0\right)>F_{2}\left(x_{0}-0\right)$ or $F_{1}\left(x_{0}+0\right)>F_{2}\left(x_{0}+0\right)$. In the first case there is an interval satisfying the conditions to the left of $x_{0}$, in the second case one to the right.
c If we remember our understanding, at the end of $\S 6.19$, about the definition of $\phi(x)$ outside $(A, B)$.
6.21. Remarks on the characteristic properties. We have still to show that the properties [1]-[3] are sufficient to characterise the means $\mathfrak{M}_{\phi}$. We insert first some general remarks concerning the 'significance' of the properties.
(i) [1] asserts that 'if all the elements of a set have the same value, then their mean has that value'.
(ii) [2] asserts that ' $\mathfrak{M}[F]$ is a strictly monotonic functional of $F^{\prime}$. It would not be sufficient to assert that (under the conditions stated) $\mathfrak{M}\left[F_{1}\right] \leqq \mathfrak{M}\left[F_{2}\right]$, i.e. that $\mathfrak{M}[F]$ is a monotonic functional'.
Let us consider some examples.
(a) The arithmetic mean

$$
\mathfrak{A}(a, q)=\Sigma q a=\int x d F=\mathfrak{A}[F]
$$

is a strictly monotonic functional of $F$. In this case $\phi(x)=x$.
(b) We may define ' $\operatorname{Max} a$ ' as 'the lower bound of the values of $x$ for which $F^{\prime}(x)=l^{\prime}$ ( $F$ being any finite distribution function with jumps at the points $a$ ). Then Max $a=\mu[F]$ is a functional of $F$ which is plainly monotonic: if $F_{1} \geqq F_{2}$ for all $x$, then $\mu\left[F_{1}\right] \leqq \mu\left[F_{2}\right]$. But $\mu[F]$ is not strictly monotonic: if $F_{1}$ and $F_{2}$ are defined by

$$
F_{1}=F_{2}=0(x<0) ; \quad F_{1}=\frac{1}{2}, \quad F_{2}=0(0<x<1) ; \quad F_{1}=F_{2}=1(x>1),
$$

then $\quad \mu\left[F_{1}\right]=\operatorname{Max}(0,1)=\operatorname{Max}(1,1)=\mu\left[F_{2}\right]$.
That $\mu[F]$ is not representable in the form (6.20.1) follows from the theorem itself; if it were, it would be strictly monotonic.
(c) The geometric mean $\mathfrak{G}=\mathfrak{G}(a, q)$ is a functional of $F$ which is not strictly monotonic, since, for example, the sets ( $0, a_{2}, \ldots$ ) and ( $0, b_{2}, \ldots$ ) have the same $\mathfrak{G b}$. It is representable by the formula

$$
\mathfrak{G}=\exp \left(\int_{0}^{\infty} \log x d F^{\prime}(x)\right) .
$$

This is of the form (6.20.1), with $\phi(x)=\log x$ for $x>0$; but $\mathfrak{f}$ is not represented in the manner prescribed by the theorem, since $\log x \rightarrow-\infty$ when $x \rightarrow 0$.
(iii) If we use [3] twice, the second time with $F^{*}, G, G^{*}, 1-t$ in place of $G, F, F^{*}, t$, we see that

$$
\begin{equation*}
\mathfrak{M}[t F+(1-t) G]=\mathfrak{M}\left[\left[t F^{*}+(1-t) G^{*}\right]\right. \tag{6.21.1}
\end{equation*}
$$

whenever $\mathfrak{M}[F]=\mathfrak{M}\left[F^{*}\right]$ and $\mathfrak{M}[G]=\mathfrak{M}\left[G^{*}\right]$. In other words
(a) $\mathfrak{M}[t F+(1-t) G]$ is determined uniquely by $\mathfrak{M}[F], \mathfrak{M}[G]$ and $t$.

More generally
(6.21.2) $\quad \mathfrak{M}\left[\Sigma q_{\nu} F_{\nu}\right]=\mathfrak{M}\left[\Sigma q_{\nu} F_{\nu}{ }^{*}\right]$
if $\mathfrak{M}\left[F_{\nu}\right]=\mathfrak{M}\left[F_{\nu}{ }^{*}\right]$ and $\Sigma q_{\nu}=1$.
A functional $\mathfrak{F}[F]$ is said to be linear if

$$
\mathfrak{F}[t F+u G]=t \mathfrak{F}[F]+u \mathfrak{F}[G]:
$$

in this case it has certainly property (a). If $\mathfrak{M}[F]$ satisfies (a) or [3], of which $(a)$ is a consequence, we may call $\mathfrak{M}[F]$ quasilinear. If also we agree to describe the property [1], as is natural, by consistency, we may state Theorem 215 shortly as follows: the most general consistent, strictly increasing, and quasi-linear functional of $F$ is that defined by (6.20.1).
6.22. Completion of the proof of Theorem 215. The functions $E_{A}(x), E_{B}(x)$, and $(1-t) E_{A}(x)+t E_{B}(x)$, where $0<t<1$, belong to $\mathfrak{D}(A, B)$. $^{\text {a }}$ We write

$$
\psi(t)=\mathfrak{M}\left[(1-t) E_{A}+t E_{B}\right],
$$

so that $\quad \psi(0)=\mathfrak{M}\left[E_{A}\right]=A, \quad \psi(1)=\mathfrak{M}\left[E_{B}\right]=B$.
Let us assume provisionally that $\psi(t)$ is strictly increasing and continuous. Then $\psi(t)$ has an inverse

$$
\phi(u)=\psi^{-1}(u)
$$

which is also continuous and increases strictly from 0 to $l$ when $u$ increases from $A$ to $B$. If

$$
u=\psi(t), \quad t=\phi(u),
$$

then $\quad \mathfrak{M}\left[E_{u}\right]=u=\psi(t)=\mathfrak{M}\left[(1-\phi(u)) E_{A}+\phi(u) E_{B}\right]$.
Hence, using [3] in the extended form (6.21.2), and the expression (6.19.1) for any finite distribution function $F$, we obtain

$$
\begin{aligned}
\mathfrak{M}[F] & =\mathfrak{M}\left[\Sigma q E_{a}\right] \\
& =\mathfrak{M}\left[\Sigma q\left\{(1-\phi(a)) E_{A}+\phi(a) E_{B}\right\}\right] \\
& =\mathfrak{M}\left[(1-\Sigma q \phi(a)) E_{A}+(\Sigma q \phi(a)) E_{B}\right] \\
& =\psi(\Sigma q \phi(a))=\phi^{-1}(\Sigma q \phi(a)),
\end{aligned}
$$

the result of the theorem.
It should be observed that here $\phi(A)=0, \phi(B)=1$. When a $\phi$ has been found, it may (after Theorem 83) be replaced by any $\alpha \phi+\beta$.
${ }^{\text {a }} E_{A}$ and $E_{B}$ are extreme cases of functions of $\mathfrak{D}(A, B)$ : if $F$ belongs to $\mathfrak{D}(A, B)$, then $E_{A} \geqq F \geqq E_{B}$ for all $x$.

It has still to be proved that $\psi(t)$ is strictly increasing and continuous.
(1) If

$$
0 \leqq t_{1}<t_{2} \leqq 1
$$

then

$$
\left(1-t_{1}\right) E_{A}+t_{1} E_{B} \geqq\left(1-t_{2}\right) E_{A}+t_{2} E_{B}
$$

for all $x$, with inequality for some $x$. Hence, by [2],

$$
\psi\left(t_{1}\right)=\mathfrak{M}\left[\left(1-t_{1}\right) E_{A}+t_{1} E_{B}\right]<\mathfrak{M}\left[\left(1-t_{2}\right) E_{A}+t_{2} E_{B}\right]=\psi\left(t_{2}\right) .
$$

(2) Suppose, if possible, that $\psi(t)$ has a discontinuity on the right at $t_{0}$, where $0 \leqq t_{0}<1$. Then we can find a $\xi$ such that

$$
\psi\left(t_{0}\right)<\xi<\psi\left(t_{0}+\epsilon\right)
$$

for arbitrarily small $\epsilon$, and

$$
E_{\psi\left(t_{0}\right)} \geqq E_{\xi} \geqq E_{\psi\left(t_{0}+\epsilon\right)}
$$

for all $x$, with inequality for some $x$. Hence, by [2],

$$
\begin{equation*}
\mathfrak{M}\left[\frac{1}{2} E_{\psi\left(t_{0}\right)}+\frac{1}{2} E_{\psi(t)}\right]<\mathfrak{M}\left[\frac{1}{2} E_{\xi}+\frac{1}{2} E_{\psi(t)}\right]<\mathfrak{M}\left[\frac{1}{2} E_{\psi\left(t_{0}+\epsilon\right)}+\frac{1}{2} E_{\psi(t)}\right] \tag{6.22.1}
\end{equation*}
$$

for any $t$ of $(0,1)$. But if $s$ and $t$ lie in $(0,1)$, then, by [1],

$$
\psi(s)=\mathfrak{M}\left[E_{\psi(s)}\right]=\mathfrak{M}\left[(1-s) E_{A}+s E_{B}\right]
$$

and similarly for $t$; and, by [3],

$$
\begin{aligned}
\mathfrak{M}\left[\frac{1}{2} E_{\psi(s)}+\frac{1}{2} E_{\psi(t)}\right] & =\mathfrak{M}\left[\frac{(1-s) E_{A}+s E_{B}}{2}+\frac{(1-t) E_{A}+t E_{B}}{2}\right] \\
& =\mathfrak{M}\left[\left(1-\frac{s+t}{2}\right) E_{A}+\frac{s+t}{2} E_{B}\right]=\psi\left(\frac{s+t}{2}\right) .
\end{aligned}
$$

Combining this with (6.22.1), we see that

$$
\psi\left(\frac{t_{0}+t}{2}\right), \quad \psi\left(\frac{t_{0}+t+\epsilon}{2}\right)
$$

are separated by a number, viz. $\mathfrak{M}\left[\frac{1}{2} E_{\xi}+\frac{1}{2} E_{\psi(t)}\right]$, which is independent of $\epsilon$; and so, making $\epsilon \rightarrow 0$,

$$
\psi\left(\frac{t_{0}+t}{2}\right)<\psi\left(\frac{t_{0}+t}{2}+0\right) .
$$

Hence $\psi$ has a discontinuity at $\frac{1}{2}\left(t_{0}+t\right)$, for all $t$ of an interval; and this is impossible, because the discontinuities of a monotonic function are at most enumerable.

It follows that $\psi(t)$ has no right-hand discontinuity. Similarly, it has no left-hand discontinuity. It is therefore continuous, and this completes the proof.

We have confined our attention to finite distribution functions, so that all the functions $F$ which have been considered are step-functions, and the means are the means of Ch. III. There is a similar theorem in which both hypothesis and conclusion are stronger in that they apply to a class of functions more extensive than $\mathfrak{D}(A, B)$. Let us denote by $\mathfrak{D}^{*}(A, B)$ the class of functions which have the properties (ii) and (iii) of §6.19 and also satisfy (6.19.4). We can then prove a theorem which differs from Theorem 215 only in the substitution of $\mathfrak{D}^{*}$ for $\mathfrak{D}$. The proof is very much the same, but is slightly more elaborate in its final stages. See de Finetti (1).

## MISCELLANEOUS THEOREMS AND EXAMPLES

216. 'Velocity averaged by time is less than velocity averaged by distance.'
[This is $\left(\int \frac{d s}{d t} d t\right)^{2}<\int d t \int \frac{d s}{d t} d s=\int d t \int\left(\frac{d s}{d t}\right)^{2} d t$, a case of Theorem 181.]
217. If the kinetic energy of a mass $M$ of moving homogeneous incompressible fluid is $E$, and the average velocity of its particles is $V$, then $E>\frac{1}{2} M V^{2}$, unless all particles have the same velocity.
[If $\rho$ is the density, $v$ the velocity of an element $d S=d x d y d z$, then

$$
M=\rho \int d S, \quad V \int d S=\int v d S, E=\frac{1}{2} \rho \int v^{2} d S,
$$

and the result follows from Theorem 181 (for triple integrals).]
218. A unit electric current passes through a closed plane circuit enclosing an area $A$, and exerts a force $F$ on a unit magnetic pole $P$ in the plane of and interior to the circuit. Then

$$
2 A F^{2}>(2 \pi)^{3}
$$

unless the circuit is a circle whose centre is $P$.
[Suppose, for simplicity, that the circuit is 'star shaped' with respect to $P$ (i.e. that every point of the line from $P$ to any point of the circuit lies inside the circuit). Then, using polar coordinates $r, \theta$ about $P$, and integrals from 0 to $2 \pi$,

$$
2 \pi=\int d \theta=\int\left(\frac{1}{r}\right)^{\frac{2}{3}}\left(r^{2}\right)^{\frac{1}{3}} d \theta<\left(\int \frac{d \theta}{r}\right)^{\frac{2}{3}}\left(\int r^{2} d \theta\right)^{\frac{1}{3}}=F^{\frac{2}{3}}(2 A)^{\frac{1}{3}},
$$

unless $r$ is constant.]
219. If $f_{\nu}(x, y)$ and $g_{\nu}(x, y)$ are two (finite or infinite) sets of functions of $x$ and $y$, then $\left(\Sigma \iint f g d x d y\right)^{2}<\Sigma \iint f^{2} d x d y \Sigma \iint g^{2} d x d y$, unless there are two constants $a$ and $b$, not both zero, such that

$$
a f_{\nu}(x, y) \equiv b g_{\nu}(x, y),
$$

for every $\nu$.
[From Theorems 7 (for infinite series) and 181 (for double integrals), or directly, by the second method of $\S 2.4$. The theorem illustrates the following principle. The inequality
where $u$ and $v$ are functions of three integral variables $m, n$ and $p$, does not differ materially from the ordinary form of Cauchy's inequality; but we can derive materially different inequalities from (i) by replacing different selections of the signs of summation by signs of integration.]
220. Suppose that the $a$ are positive, and that $q_{r}$ is defined by

$$
\frac{1}{\left(1-a_{1} x\right)\left(1-a_{2} x\right) \ldots\left(1-a_{n} x\right)}=1+n q_{1} x+\ldots+\binom{n+r-1}{r} q_{r} x^{r}+\ldots
$$

Then $\quad q_{r}^{2}<q_{r-1} q_{r+1} \quad(r=1,2, \ldots)$
unless all the $a$ are equal.
221.

$$
q_{1}<q_{2}{ }^{\frac{1}{2}}<q_{3}{ }^{\frac{1}{3}}<\ldots
$$

unless all the $a$ are equal.
[Theorems 220 and 221 were communicated to us by Prof. I. Schur. The $q$ are means of homogeneous products of the $a$, like the $p$ of $\S 2.22$, but now the $a$ in a product are not necessarily different. In particular $q_{1}=p_{1}$.

Theorem 221 follows from Theorem 220 as Theorem 52 followed from Theorem 51. To prove Theorem 220 we observe that

$$
\begin{equation*}
q_{r}=(n-1)!\iint \ldots \int\left(a_{1} x_{1}+a_{2} x_{2}+\ldots+a_{n} x_{n}\right)^{r} d x_{1} \ldots d x_{n-1} \tag{i}
\end{equation*}
$$

where $x_{n}=1-x_{1}-x_{2}-\ldots-x_{n-1}$ and the domain of integration is defined by $x_{1}>0, \ldots, x_{n-1}>0, x_{n}>0$. We obtain Theorem 220 by applying Theorem 181 (for multiple integrals) to (i).

The formula (i) leads to a more complete theorem. If the $a$ are real (but not necessarily positive) then the quadratic form $\Sigma q_{r+s} y_{r} y_{s}$ is strictly positive; and if the $a$ are positive, then the form $\Sigma q_{r+s+1} y_{r} y_{s}$ is strictly positive; except (in both cases) when all the $a$ are equal.]
222. If $p>1, f$ is $L^{p}$ in ( $\left.0, a\right)$, and
then

$$
\begin{gathered}
F(x)=\int_{0}^{x} f(t) d t, \\
F(x)=o\left(x^{1 / p}\right)
\end{gathered}
$$

for small $x$.
[By Theorem 189,

$$
F^{p} \leqq \int_{0}^{x} f^{p} d t\left(\int_{0}^{x} d t\right)^{p-1}=x^{p-1} \int_{0}^{x} f^{p} d t
$$

and the second factor tends to 0.]
223. If $p>1$ and $f$ is $L^{p}$ in $(0, \infty)$, then $F(x)=o\left(x^{1 / p^{\prime}}\right)$ both for small and for large $x$.
[For small $x$, by Theorem 222. To prove the result for large $x$, choose $X$ so that

$$
\int_{X}^{\infty} f^{p} d x<\epsilon^{p}
$$

and suppose $x>X$. Then

$$
\begin{gathered}
(F(x)-F(X))^{p}=\left(\int_{X}^{x} f d t\right)^{p} \leqq(x-X)^{p-1} \int_{X}^{x} f^{p} d t<\epsilon^{p} x^{p-1}, \\
F(x)<F(X)+\epsilon x^{1 / p^{\prime}}<2 \epsilon x^{1 / p^{\prime}}
\end{gathered}
$$

for sufficiently large $x$.]
224. If $y$ is an integral except perhaps at $x=0$ and $x y^{\prime 2}$ is integrable in ( $0, a$ ), then

$$
y=o\left\{\left(\log \frac{1}{x}\right)^{\frac{1}{2}}\right\}
$$

for small $x$.

$$
\left[\int_{e}^{a} y^{\prime} d x \leqq\left(\int_{\epsilon}^{a} \frac{d x}{x} \int_{\epsilon}^{a} x y^{\prime 2} d x\right)^{\frac{1}{2}} \cdot\right]
$$

225. If $y$ is an integral except perhaps at 0 and 1 , and $x(1-x) y^{\prime 2}$ is integrable in ( 0,1 ), then $y$ is $L^{2}$ and

$$
0 \leqq \int_{0}^{1} y^{2} d x-\left(\int_{0}^{1} y d x\right)^{2} \leqq \frac{1}{2} \int_{0}^{1} x(1-x) y^{\prime 2} d x
$$

[That $y$ is $L^{2}$ follows from Theorem 224. The first inequality is included in Theorem 181. For the second, we have

$$
\begin{aligned}
& \int_{0}^{1} y^{2} d x-\left(\int_{0}^{1} y d x\right)^{2}=\frac{1}{2} \int_{0}^{1} \int_{0}^{1}\{y(u)-y(v)\}^{2} d u d v \\
& \quad=\int_{0}^{1} d u \int_{u}^{1} d v\left(\int_{u}^{v} y^{\prime}(t) d t\right)^{2} \leqq \int_{0}^{1} d u \int_{u}^{1}(v-u) d v \int_{u}^{v}\left(y^{\prime}(t)\right)^{2} d t \\
& \quad=\int_{0}^{1}\left(y^{\prime}(t)\right)^{2} d t \int_{0}^{t} d u \int_{t}^{1}(v-u) d v=\frac{1}{2} \int_{0}^{1} t(1-t) y^{\prime 2} d t
\end{aligned}
$$

Of the two inequalities, the first can reduce to an equality only if $y$ is constant, the second only if $y$ is linear.]
226. If $m>1, n>-1$, and $f$ is positive and an integral, then

$$
\begin{equation*}
\int_{0}^{\infty} x^{n} f^{m} d x \leqq \frac{m}{n+1}\left\{\int_{0}^{\infty} x^{\frac{m(n+1)}{m-1}} f^{m} d x\right\}^{\frac{m-1}{m}}\left(\int_{0}^{\infty}\left|f^{\prime}\right|^{m} d x\right)^{\frac{1}{m}} \tag{i}
\end{equation*}
$$

with equality only when $f=B \exp \left\{-C x^{(m+n) /(m-1)}\right\}$, where $B \geqq 0, C>0$.
In particular

$$
\begin{equation*}
\int_{0}^{\infty} f^{2} d x<2\left(\int_{0}^{\infty} x^{2} f^{2} d x\right)^{\frac{1}{2}}\left(\int_{0}^{\infty} f^{\prime 2} d x\right)^{\frac{1}{2}}, \tag{ii}
\end{equation*}
$$

unless $f=B e^{-C x^{2}}$; and this inequality holds whether $f$ be positive or not, and also for the range $(-\infty, \infty)$.
[The most interesting case is (ii), which is due to Weyl (1,345), and is useful in quantum-mechanics.

Assume that the integrals on the right-hand side of (i) are finite. Since $f$ is continuous, and $n<m(n+1) /(m-1)$, that on the left-hand side is also finite. Hence

$$
{\underset{x}{\lim }}_{\rightarrow \infty} x^{n+1} f^{m}=0,
$$

and so, integrating by parts over ( $0, x_{k}$ ), where $\left(x_{k}\right)$ is an appropriate sequence which tends to infinity with $k$,

$$
\int_{0}^{\infty} x^{n} f^{m} d x=-\frac{m}{n+1} \lim _{k \rightarrow \infty} \int_{0}^{x_{k}} x^{n+1} f^{m-1} f^{\prime} d x
$$

But, by Theorem 189,

$$
\int_{0}^{\infty} x^{n+1} f^{m-1}\left|f^{\prime}\right| d x<\left\{\int_{0}^{\infty} x^{\frac{m(n+1)}{m-1}} f^{m} d x\right\}^{\frac{m-1}{m}}\left(\int_{0}^{\infty}\left|f^{\prime}\right|^{m} d x\right)^{\frac{1}{m}},
$$

unless $f^{\prime} \leqq 0$ and $f^{\prime}$ and $x^{(n+1) /(m-1)} f$ are effectively proportional. This hypothesis leads to the form of $f$ stated.]
227. If $\phi$ increases,

$$
\left(\int f g d \phi\right)^{2}<\int f^{2} d \phi \int g^{2} d \phi,
$$

unless $f$ and $g$ are effectively proportional (in the sense of $\S 6.17$ ).
[Included in Theorem 210; wanted in Theorem 228.]
228. If $a \geqq 0, b \geqq 0, a \neq b$, and $\phi$ is non-negative and decreasing, then

$$
\left(\int_{0}^{\infty} x^{a+b} \phi d x\right)^{2}<\left\{1-\left(\frac{a-b}{a+b+1}\right)^{2}\right\} \int_{0}^{\infty} x^{2 a} \phi d x \int_{0}^{\infty} x^{2 b} \phi d x
$$

unless $\phi=C$, where $C>0$, in $(0, \xi)$, and $\phi=0$ in $(\xi, \infty)$.
[ $\xi$ may be 0 . The inequality is stronger than that resulting from a direct application of Theorem 181. It follows from Theorem 227 if we reduce the integrals to the form considered there by partial integration. The case $a=0, b=2$ was mentioned by Gauss in connection with the Theory of Errors: see Gáuss (1, iv, 12) and Pólya and Szegö (1, iI, 114, 318).]
229. If $a \geqq 0, b \geqq 0, a \neq 1$, and $\phi$ is non-negative and increasing, then

$$
\left(\int_{0}^{1} x^{a+b} \phi d x\right)^{2}>\left\{1-\left(\frac{a-b}{a+b+1}\right)^{2}\right\} \int_{0}^{1} x^{2 a} \phi d x \int_{0}^{1} x^{2 b} \phi d x,
$$

unless $\phi=C$.
[See Pólya and Szegö (1, 1, 57, 214). In this case Theorem 181 gives a reversed inequality, with the factor 1 on the right-hand side.]
230. If $\quad 0<a \leqq f \leqq A<\infty, \quad 0<b \leqq g \leqq B<\infty$,
then

$$
\int f^{2} d x \int g^{2} d x \leqq\left[\frac{1}{2}\left\{\sqrt{ }\left(\frac{A B}{a b}\right)+\sqrt{ }\left(\frac{a b}{A B}\right)\right\} \int f g d x\right]^{2}
$$

[Analogue of Theorem 71: see Pólya and Szegö (1, r, 57, 214).]
231. If we consider the closed or open intervals (in general four in number) with end-points $-a, b$, where $a \geqq 0, b \geqq 0$, and suppose each of $a$ and $b$ zero, positive and finite, or infinite, we obtain in all 34 types of intervals $I$. Assign to each interval $I$ a function $f(x)$ defined for $0<x<1$ and such that $\log \mathfrak{M}_{r}(f)$, where $\mathfrak{M}_{r}(f)$ is formed for the interval $(0,1)$ and with $q=1$, is finite just for the values of $r$ in $I$.
[Examples:
$I$ is $-a<r \leqq b$;

$$
I \text { is }-\infty \leqq r \leqq \infty \text {; }
$$

$$
I \text { is the single point } 0 \text {; }
$$

$$
\begin{aligned}
& f(x)=x^{1 / a}(1-x)^{-1 / b}\left(\log \frac{2}{1-x}\right)^{-2 / b}: \\
& f(x)=1+x^{2}: \\
& f(x)=\exp \left(-x^{-\frac{1}{2}}+(1-x)^{-\frac{1}{2}}\right): \\
& f(x)=\exp \left(-x^{-1}+(1-x)^{-1}\right) .
\end{aligned}
$$

This contains part of the proof of what is stated near the end of § 6.11.]
232. Geometrical interpretation of Minkowski's inequality. Suppose that a point in functional space is defined as a function of $L^{2}$, two functions defining the same point if and only if their difference is nul; and that the distance between two points $f$ and $g$ is defined by

$$
\delta(f, g)=\sqrt{ }\left\{\int(f-g)^{2} d x\right\} .
$$

Then (i) the distance between two distinct points is positive; and (ii)

$$
\delta(f, h) \leqq \delta(f, g)+\delta(g, h) .
$$

[If we define distance by

$$
\delta(f, g)=\left(\int|f-g|^{r} d x\right)^{1 / r} \quad(r \geqq 1)
$$

we obtain similar results in 'functional space $L^{r}$ '.]
233. The shortest distance between two given points in Euclidean space is the straight line.
[A curve in space is given by

$$
x=x(t), \quad y=y(t), \quad z=z(t) .
$$

We may suppose that $t$ increases from 0 to 1 on the arc in question. If we assume that $x, y, z$ are integrals of functions of $L^{2}$, then the length $l$ is given by
$l^{2}=\left[\int\left(x^{\prime 2}+y^{\prime 2}+z^{\prime 2}\right)^{\frac{1}{2}} d t\right]^{2}=\mathfrak{M}_{\frac{1}{2}}\left(x^{\prime 2}+y^{\prime 2}+z^{\prime 2}\right) \geqq \mathfrak{M}_{\frac{1}{2}}\left(x^{\prime 2}\right)+\mathfrak{M}_{\frac{1}{2}}\left(y^{\prime 2}\right)+\mathfrak{M}_{\frac{1}{2}}\left(z^{\prime 2}\right)$, by Theorem 198; and this is not less than

$$
\left(\int x^{\prime} d t\right)^{2}+\left(\int y^{\prime} d t\right)^{2}+\left(\int z^{\prime} d t\right)^{2}=\left(x_{1}-x_{0}\right)^{2}+\left(y_{1}-y_{0}\right)^{2}+\left(z_{1}-z_{0}\right)^{2} .
$$

If there is equality, $A x^{\prime} \equiv B y^{\prime} \equiv C z^{\prime}$, and the curve is a straight line.]
234. If $0<p<1$ and

$$
\int f g d x \geqq A\left(\int g^{p^{\prime}} d x\right)^{1 / p^{\prime}}
$$

for all $g$, then $\int f^{p} d x \geqq A^{p}$.
[Compare Theorem 70. If $f>0$ for all $x$, define $g$ by $f g=f^{p}$. If $f>0$ in $E, f=0$ in $C E$, and the measure of $C E$ is finite, define $g$ by $f g=f^{p}$ in $E$ and by $g=G$ in $C E$, and proceed as in the proof of Theorem 70. If the measure of $C E$ is infinite, take (for example)
in $C E$. Then

$$
g=G e^{x^{2}}
$$

$$
\int f^{p} d x=\int f g d x \geqq A\left(\int_{E} f^{p} d x+G^{p^{\prime}} \int_{C E} e^{p^{\prime} x^{2}} d x\right)^{1 / p^{\prime}}
$$

and the result again follows when $G \rightarrow \infty$.]
235. Suppose that $f$ and $p$ are positive and that $f$ has the period $2 \pi$; that

$$
F(x)=\int_{0}^{2 \pi} f(x+t) p(t) d t / \int_{0}^{2 \pi} p(t) d t ;
$$

and that the means $\mathfrak{M}_{r}$ refer to the interval ( $0,2 \pi$ ) and a constant weightfunction. Then

$$
\mathfrak{M}_{r}(F) \geqq \mathfrak{M}_{r}(f) \quad(0 \leqq r \leqq 1), \quad \mathfrak{M}_{r}(F) \leqq \mathfrak{M}_{r}(f) \quad(r \geqq 1) .
$$

[This may be deduced from Theorem 204, or proved directly (supposing for example $r \geqq 1$ ) as follows:

$$
\begin{aligned}
\mathfrak{M}_{r}^{r}(F) & =\frac{1}{2 \pi} \int d x\left[\frac{\int f(x+t) p(t) d t}{\int p(t) d t}\right]^{r} \\
& \leqq \frac{1}{2 \pi} \int d x \frac{\int f^{r}(x+t) p(t) d t\left(\int p(t) d t\right)^{r-1}}{\left(\int p(t) d t\right)^{r}} \\
& =\frac{1}{2 \pi} \frac{\int p(t) d t \int f^{r}(x+t) d x}{\int p(t) d t}=\frac{1}{2 \pi} \int f^{r}(u) d u=\mathfrak{M}_{r}^{r}(f) .
\end{aligned}
$$

For the case $r=0$, see Pólya and Szegö (1, r, 56, 212).]
236. We say that $f(x, y, \ldots)$ and $g(x, y, \ldots)$ are similarly ordered if

$$
\left\{f\left(x_{1}, y_{1}, \ldots\right)-f\left(x_{2}, y_{2}, \ldots\right)\right\}\left\{g\left(x_{1}, y_{1}, \ldots\right)-g\left(x_{2}, y_{2}, \ldots\right)\right\} \geqq 0,
$$

oppositely ordered if $f$ and $-g$ are similarly ordered. Prove that

$$
\iint \ldots f d x d y \ldots \iint \ldots g d x d y \ldots \leqq \iint \ldots d x d y \ldots \iint \ldots f g d x d y \ldots
$$

if $f$ and $g$ are similarly ordered, while the sign is reversed if $f$ and $g$ are oppositely ordered. The integration is extended over any common part of the regions of definition of $f$ and $g$.
[Analogue of Theorem 43 (with $r=1$ ), due in substance to Tchebychef (who considers only monotonic functions of one variable).]
237. If $\phi$ and $\psi$ satisfy the conditions of Theorem 156, and

$$
\Phi(x)=\int_{0}^{x} \phi(t) d t, \quad \Psi(x)=\int_{0}^{x} \psi(t) d t,
$$

then

$$
\int f g d x \leqq \int \Phi(f) d x+\int \Psi^{\prime}(g) d x
$$

238. If $f$ and $g$ are positive, and $k$ a positive constant, and $f \log ^{+} f$ and $e^{k g}$ are integrable, then $f g$ is integrable.
[By Theorem 63, $k f g \leqq f \log ^{+} f+e^{k g-1}$.]
239. If $f$ is positive, then

$$
\int_{0}^{a} f \log \frac{1}{x} d x \leqq 2 \int_{0}^{a} f \log ^{+} f d x+\frac{4 \sqrt{ } a}{e} .
$$

[Take $g=\frac{1}{2} \log \frac{1}{x}, k=1$, in the inequality used in proving Theorem 238.]

240 If $f$ is positive and $L$ in $(0, a)$, and

$$
F(x)=\int_{0}^{x} f d t
$$

then

$$
\int_{0}^{a} f(x) \log \frac{1}{x} d x=\int_{0}^{a} \frac{F(x)}{x} d x+F(a) \log \frac{1}{a},
$$

whenever either integral is finite.
241. Suppose that $a$ is positive and finite; that $B=B(a)$ denotes generally a number depending on $a$ only; that $f(x) \geqq 0$; and that

$$
\begin{gathered}
F(x)=\int_{0}^{x} f(t) d t \\
J=\int_{0}^{a} f \log ^{+} f d x, \quad K=\int_{0}^{a} \frac{F}{x} d x
\end{gathered}
$$

Then (i) if $J$ is finite $K$ is also finite, and

$$
K<B J+B:
$$

(ii) when $f$ is a decreasing function the converse is also true; if $K$ is finite then $J$ is finite, and

$$
J<B K \log ^{+} K+B
$$

[For the last two theorems see Hardy and Littlewood (8).]
242. If $f$ is positive and $L$ in $(0, a)$ and

$$
\begin{gathered}
g=\int_{x}^{a} \frac{f(t)}{t} d t \\
\int_{0}^{a} g(x) d x=\int_{0}^{a} f(x) d x
\end{gathered}
$$

then $g$ is $L$ and
[Integrate by parts; or substitute for $g$ and change the order of integration.]
243. We define $\mathfrak{M}_{\phi}(f)$, where $\phi$ is a continuous and strictly increasing function, by

Then, in order that

$$
\begin{gathered}
\mathfrak{M}_{\phi}(f)=\phi^{-1}\left\{\int \phi(f) q d x\right\} . \\
\mathfrak{M}_{\phi}(f) \leqq \mathfrak{M}_{\psi}(f)
\end{gathered}
$$

for all $f$, it is necessary and sufficient that $\psi$ should be convex with respect to $\phi$.
244. In order that

$$
\mathfrak{M}_{\phi_{n}}^{x_{n}} \ldots \mathfrak{M}_{\phi_{1}}^{x_{1}}(f) \leqq \mathfrak{M}_{\psi_{n}}^{x_{n}} \ldots \mathfrak{M}_{\psi_{1}}^{x_{1}}(f)
$$

for all $f=f\left(x_{1}, x_{2}, \ldots, x_{n}\right)$, it is necessary and sufficient that every $\psi_{\nu}$ be convex with respect to the corresponding $\phi_{\nu}$.
245. In order that

$$
\mathfrak{M}_{p_{n}}^{x_{n}} \ldots \mathfrak{M}_{p_{1}}^{x_{1}}(f) \leqq \mathfrak{M}_{q_{n}}^{x_{\nu_{n}}} \ldots \mathfrak{M}_{q_{v_{1}}}^{x_{\nu_{1}}}(f)
$$

for all $f$, it is necessary and sufficient that (i) $q_{\nu} \geqq p_{\nu}$ and (ii) $q_{\mu} \geqq p_{\nu}$ when $\mu>\nu$ and the permutation by which $\nu_{1}, \nu_{2}, \ldots, \nu_{n}$ is derived from $1,2, \ldots, n$ involves an inversion of $\mu$ and $\nu$.
[For the last three theorems, which correspond to Theorems 92, 93 and 137, see Jessen (2, 3).]
246. Hölder's inequality may be deduced from

$$
\mathfrak{M}_{1}^{y} \mathfrak{M}_{0}^{x}(f) \leqq \mathfrak{M}_{0}^{x} \mathfrak{M}_{1}^{y}(f)
$$

(Theorem 203), by taking

$$
f(x, y)=f_{1}^{p}(y)\left(0 \leqq x<\frac{1}{p}\right), \quad f(x, y)=f_{2}{ }^{p \prime}(y)\left(\frac{1}{p} \leqq x \leqq 1\right) .
$$

[See Jessen (3).]
247. If (i) $\phi(x, t)$ is positive, continuous, and convex in $x$, for $x_{1} \leqq x \leqq x_{2}$, $t>0$; (ii) $p(t) \geqq 0$; (iii) the integral

$$
I(x)=\int_{0}^{\infty} \phi(x, t) p(t) d t
$$

is finite for $x=x_{1}$ and $x=x_{2}$; then $I(x)$ is continuous and convex for

$$
x_{1}<x<x_{2} .
$$

[That $I(x)$ is bounded and convex follows immediately from the convexity of $\phi$; that it is continuous, from Theorem 111.]
248. If $f(x)$ and $\phi(x)$ are positive and $\phi(x)$ convex for positive $x$, and

$$
I(x)=x \int_{0}^{\infty} \phi\left\{\frac{f(t)}{x}\right\} d t
$$

is finite for $x=x_{1}$ and $x=x_{2}$, then $I(x)$ is continuous and convex for

$$
x_{1}<x<x_{2} .
$$

[By Theorem 119, $x \phi(1 / x)$ and

$$
f(t) \cdot \frac{x}{f(t)} \phi\left\{\frac{f(t)}{x}\right\}
$$

are convex, and we can apply Theorem 247. More general results can be derived from Theorem 120.]
249. In order that

$$
\int_{a}^{b} \phi(g(x)) d x \leqq \int_{a}^{b} \phi(f(x)) d x
$$

should be true for every convex and continuous $\phi$, it is necessary and sufficient that

$$
\begin{gathered}
\int_{a}^{b} g(x) d x=\int_{a}^{b} f(x) d x \\
\int_{a}^{b}(g(x)-y)^{+} d x \leqq \int_{a}^{b}(f(x)-y)^{+} d x
\end{gathered}
$$

and
for all $y$.
[Here $a^{+}$means $\operatorname{Max}(a, 0)$, as in §6.1.]
250. If $f$ and $g$ are increasing functions, then an equivalent condition is
for $a \leqq \xi \leqq b$.

$$
\int_{\xi}^{b} g(x) d x \leqq \int_{\xi}^{b} f(x) d x
$$

[For the two last theorems, which embody analogues for integrals of parts of Theorem 108, see Hardy, Littlewood, and Pólya (2).]
251. If $f_{1}(t), f_{2}(t), \ldots, f_{l}(t)$ are real and integrable in ( 0,1 ), then either (i) there is a function $x(t)$ such that

$$
\int_{0}^{1} f_{1}(t) x(t) d t>0, \ldots, \int_{0}^{1} f_{l}(t) x(t) d t>0,
$$

or (ii) there are non-negative numbers $y_{1}, y_{2}, \ldots, y_{l}$, not all zero, such that

$$
y_{1} f_{1}(t)+y_{2} f_{2}(t)+\ldots+y_{l} f_{l}(t) \equiv 0
$$

252. If $f_{1}(t), f_{2}(t), \ldots, f_{m}(t)$ are real and continuous in $(0,1)$, then either (i) there are real numbers $x_{1}, x_{2}, \ldots, x_{m}$ such that

$$
x_{1} f_{1}(t)+x_{2} f_{2}(t)+\ldots+x_{m} f_{m}(t)
$$

is non-negative for all, and positive for some, $t$ of ( 0,1 ), or (ii) there is a positive and continuous function $y(t)$ such that

$$
\int_{0}^{1} f_{1}(t) y(t) d t=0, \ldots, \int_{0}^{1} f_{m}(t) y(t) d t=0
$$

[Theorems 251 and 252 are both integral analogues of an important theorem of Stiemke (1) concerning systems of linear inequalities. Suppose that

$$
a_{\lambda \mu} \quad(\lambda=1,2, \ldots, l ; \mu=1,2, \ldots, m)
$$

is a rectangular array of $l$ rows and $m$ columns, and that

$$
\begin{aligned}
L_{\lambda}(x) & =a_{\lambda 1} x_{1}+a_{\lambda 2} x_{2}+\ldots+a_{\lambda m} x_{m}, \\
M_{\mu}(x) & =a_{1 \mu} y_{1}+a_{2 \mu} y_{2}+\ldots+a_{l \mu} y_{l} ;
\end{aligned}
$$

and consider the two problems:
(i) to find a real set $(x)$ for which

$$
L_{1}(x)>0, L_{2}(x)>0, \ldots, L_{l}(x)>0
$$

(ii) to find a non-negative and non-nul set ( $y$ ) for which

$$
M_{1}(y)=0, M_{2}(y)=0, \ldots, M_{m}(y)=0
$$

Since

$$
\Sigma y L(x)=\Sigma x M(y)
$$

the two problems cannot both be soluble for the same set (a), and Stiemke's theorem asserts that one is soluble whatever the set ( $a$ ).

Theorems 251 and 252 state analogues of Stiemke's theorem in which the $m$ columns or $l$ rows are replaced by a continuous infinity of columns or rows. These theorems, and further references to the theory of systems of linear inequalities, which we have excluded from our programme only on account of its algebraical and geometrical preliminaries, will be found in Haar (1) and Dines (1).]

## CHAPTER VII

## SOME APPLICATIONS OF THE CALCULUS OF VARIATIONS

7.1. Some general remarks. The 'simplest problem of the Calculus of Variations' is that of determining a maximum or minimum value of

$$
J(y)=\int_{x_{0}}^{x_{1}} F\left(x, y, y^{\prime}\right) d x
$$

for all functions $y=y(x)$ for which
(1) $y_{0}=y\left(x_{0}\right), y_{1}=y\left(x_{1}\right)$ are given,
(2) $y^{\prime}$ is continuous.

Let us denote this class of functions by $\Omega$. Then our object is to find a function

$$
y=Y(x)
$$

of $\mathfrak{\Omega}$, such that either

$$
J(y)<\int_{x_{0}}^{x_{1}} F\left(x, Y, Y^{\prime}\right) d x=J(Y),
$$

or $J(y)>J(Y)$, for all $y$ of $\mathfrak{\Re}$ other than $Y$. The general theory tells us that, if such a function $Y$ exists, it must satisfy 'Euler's equation'
(E)

$$
\frac{\partial F}{\partial y}=\frac{d}{d x}\left(\frac{\partial F}{\partial y^{\prime}}\right) .
$$

Let us consider some simple examples.
(i) Suppose that $\quad J(y)=\int_{0}^{1} y^{\prime 2} d x$
and $y_{0}=0, y_{1}=1$. Then ( E ) is $y^{\prime \prime}=0$, and the only solution satisfying the conditions is $y=x$. It is easy to verify that $Y=x$ does in fact give a minimum for $J(y)$. For $J(Y)=1$ and

$$
\mathrm{l}=\left(\int_{0}^{1} y^{\prime} d x\right)^{2}<\int_{0}^{1} d x \int_{0}^{1} y^{\prime 2} d x=J(y)
$$

by Theorem 181, unless $y^{\prime} \equiv 1, y=x$; so that $J>1$ for all $y$ other than $Y$.

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We have in fact proved more than the problem as stated demands, since the proof is valid whenever $y$ is an integral. This last hypothesis is however essential, since there are functions $y$ for which

$$
\begin{equation*}
y(0)=0, \quad y(1)=1, \quad y^{\prime} \equiv 0, \quad J(y)=0 . \tag{7.1.1}
\end{equation*}
$$

In order that $y$ should be an integral, i.e. in order that there should be an integrable function $f(x)$ such that

$$
y(x)=\int_{a}^{x} f(u) d u
$$

it is necessary and sufficient that $y(x)$ should be 'absolutely continuous'. It is necessary, but not sufficient, that $y(x)$ should have bounded variation. In particular it is not sufficient that $y(x)$ should be monotonic; there are increasing functions $y$ which satisfy (7.1.1).

If $y$ is the integral of $f$, then $y^{\prime} \equiv f$; an integral is the integral of its derivative. All this is expounded in detail in books on the theory of functions of a real variable ${ }^{\text {a }}$. The main theorem needed in this chapter is the theorem of integration by parts, stated in $\S 6.3(a)$.

These remarks lead us to lay down the following convention. Throughout this chapter it will be assumed that, whenever $y$ and $y^{\prime}$ occur in an enunciation or a proof, $y$ is an integral (and so the integral of $y^{\prime}$ ). A similar assumption will be made about $y^{\prime}$ and $y^{\prime \prime}$ (if $y^{\prime \prime}$ occurs in the problem); and the assumption naturally applies also to letters other than $y$. Without this assumption, all the problems of this chapter would lose their significance.
(ii) ${ }^{\text {b }}$ Suppose that

$$
J(y)=\int_{0}^{1}\left(y^{\prime 2}+y^{\prime 3}\right) d x
$$

and $y_{0}=y_{1}=0$. The only solution of (E) satisfying the conditions is $y=0$. If $Y=0, J(Y)=0$, but $Y$ does not give a maximum or minimum of $J(y)$. It is in fact easy to construct a $y$ of $\Omega$ for which

[^44]$J(y)$ is as large, positively or negatively, as we please. Thus, if $f(x)$ is any function for which $f(0)=f(1)=0$ and
$$
\int_{0}^{1} f^{\prime 3} d x>0
$$
and $y=C f(x)$, then $J(y)$ is large when $C$ is large and has the sign of $C$.
(iii) Suppose that
$$
J(y)=\int_{0}^{1} x y^{\prime 2} d x
$$
$y_{0}=0, y_{1}=1$. Here $J(y)>0$ for all $y$ of $\mathfrak{\Re}$; but $y=x^{m}$ gives $J=\frac{1}{2} m$, so that there are $y$ of $\mathfrak{\Re}$ for which $J(y)$ is as small as we please. For $y$ of $\mathfrak{K}, J(y)$ has an unattained lower bound 0 . The same is true for classes of $y$ more general than $\Re$ (for example, the class of integrals). On the other hand, $J(y)$ attains its bound 0 for the function mentioned under (i) above, and also for the discontinuous function which is 0 for $x=0$ and 1 for $x>0$.
7.2. Object of the present chapter. The Calculus of Variations might be expected to provide a very powerful weapon for the proof of integral inequalities. There are however hardly any instances of its application to inequalities of the types important in general analysis. This may be explained on two grounds. In the first place, the Calculus of Variations is concerned avowedly with attained maxima or minima, while many of the most important integral inequalities assert unattained upper or lower bounds. Secondly, the 'continuity' hypotheses of the classical theory are very restrictive. It is often more troublesome to extend an inequality, proved by variational methods for a special class of functions, to the most general classes for which the inequality is required, than to construct a direct proof of the full result. For these reasons the Calculus of Variations has been almost ignored in this chapter of analysis.

The ideas of the Calculus are however often very useful, and we apply them here to a number of special inequalities. When, as in example (i) above, or Theorems 254 and 256 below, the bound asserted by the inequality is attained, and attained by an ex-
tremal, that is to say by a solution of Euler's equation, these ideas are obviously relevant, and the result may well be one which it would be difficult to obtain in any other way. We shall however find that they are sometimes effective even when the bound is unattained and the final result lies outside the scope of the theory.

Our arguments will not demand any detailed knowledge of the theory; except in $\S 7.8$, we shall require only its simplest formal ideas ${ }^{\text {a }}$.
7.3. Example of an inequality corresponding to an unattained extremum. As a first example of the use of variational methods, we select a special case of a theorem which was first proved in an entirely different manner, and to which we shall return in § 9.8.
253. If $y^{\prime}$ belongs to $L^{2}(0, \infty), y_{0}=0$, and $y$ is notalwayszero, then

$$
J(y)=\int_{0}^{\infty}\left(4 y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x>0
$$

It is necessary for our present purpose to consider the more general integral

$$
\begin{equation*}
J(y)=\int_{0}^{\infty}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x \quad(\mu \geqq 4) \tag{7.3.1}
\end{equation*}
$$

Euler's equation is

$$
x^{2} y^{\prime \prime}+\lambda y=0 \quad\left(\lambda=1 / \mu \leqq \frac{1}{4}\right)
$$

Its solution is

$$
y=A x^{m}+B x^{n}
$$

where

$$
m=\frac{1}{2}+\sqrt{ }\left(\frac{1}{4}-\lambda\right), \quad n=\frac{1}{2}-\sqrt{ }\left(\frac{1}{4}-\lambda\right)
$$

if $\mu>4$, and is

$$
y=x^{\frac{1}{2}}(A+B \log x)
$$

if $\mu=4$. In neither case is there a solution (other than $y=0$ ) for which $y^{\prime}$ is $L^{2}$.

For this reason it is necessary to modify the problem before we attempt to apply variational ideas. We consider

$$
J(y)=\int_{0}^{1}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x
$$

[^45]with $y_{0}=0, y_{1}=1, \mu>4$. There is then one ${ }^{\text {a }}$ extremal satisfying the conditions, viz.
\[

$$
\begin{equation*}
y=Y=x^{m}=x^{\frac{1}{2}+a} \tag{7.3.2}
\end{equation*}
$$

\]

where

$$
\begin{equation*}
a=\sqrt{\left(\frac{1}{4}-\frac{1}{\mu}\right), \quad \mu=\frac{1}{\frac{1}{4}-a^{2}} . . . . ~} \tag{7.3.3}
\end{equation*}
$$

A simple calculation gives

$$
\begin{equation*}
J(Y)=\frac{2}{1-2 a} \tag{7.3.4}
\end{equation*}
$$

and this suggests the following theorem.
254. ${ }^{\text {b }}$ If $\mu>4, y(0)=0, y(1)=1$, and $y^{\prime}$ is $L^{2}$, then

$$
\begin{equation*}
J(y)=\int_{0}^{1}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x \geqq \frac{2}{1-2 a} \tag{7.3.5}
\end{equation*}
$$

where $a$ is defined by (7.3.3). The only case of equality is that defined by (7.3.2).
7.4. First proof of Theorem 254. We give two proofs of Theorem 254. The first demands no knowledge of the Calculus of Variations, though the transformations which we use are suggested by our knowledge of the form of the extremal $Y$.

If

$$
\begin{equation*}
y=x^{\frac{1}{2}+a}+\eta=Y+\eta \tag{7.4.1}
\end{equation*}
$$

then

$$
\begin{equation*}
J(y)=J(Y)+J(\eta)+K(Y, \eta) \tag{7.4.2}
\end{equation*}
$$

where

$$
K(Y, \eta)=2 \int_{0}^{1}\left(\mu Y^{\prime} \eta^{\prime}-\frac{Y \eta}{x^{2}}\right) d x
$$

Since $Y^{\prime}$ and so $\eta^{\prime}$ are $L^{2}, \eta=o\left(x^{\frac{1}{2}}\right)$ for small $x ;$ c and so

$$
\begin{aligned}
K & =2 \lim _{\delta \rightarrow 0}\left(-\int_{\delta}^{1} \frac{Y \eta}{x^{2}} d x+\left[\mu Y^{\prime} \eta\right]_{\delta}^{1}-\int_{\delta}^{1} \mu Y^{\prime \prime} \eta d x\right) \\
& =-2 \lim _{\delta \rightarrow 0} \int_{\delta}^{1} \eta\left(\mu Y^{\prime \prime}+\frac{Y}{x^{2}}\right) d x=0 .
\end{aligned}
$$

a The extremal $\quad y=\lambda x^{\frac{1}{2}}+a+(1-\lambda) x^{\frac{1}{2}-a}$
gives $y_{0}=0, y_{1}=1$ for any $\lambda$; but $y^{\prime}$ is not $L^{2}$, and $J(y)$ diverges, unless $\lambda=1$.
${ }^{b}$ For this and some later theorems in this chapter see Hardy and Littlewood (10).
c Theorem 222.

Hence (7.4.2) gives

$$
\begin{equation*}
J(y)=\frac{2}{1-2 a}+J(\eta) \tag{7.4.3}
\end{equation*}
$$

and it is sufficient to prove that

$$
J(\eta)>0
$$

Here $\eta^{\prime}$ is $L^{2}$, and $\eta$ is not nul but vanishes at the ends of the interval.

We now write $\quad \eta=Y \zeta$.
Then
(7.4.4)

$$
\text { 4) } \begin{aligned}
& J_{\delta}(\eta)=\int_{\delta}^{1}\left(\mu \eta^{\prime 2}-\frac{\eta^{2}}{x^{2}}\right) d x \\
= & \dot{\mu} \int_{\delta}^{1}\left(Y \zeta^{\prime}+Y^{\prime} \zeta\right)^{2} d x-\int_{\delta}^{1} \frac{Y^{2} \zeta^{2}}{x^{2}} d x \\
= & \mu \int_{\delta}^{1} \dot{Y^{2} \zeta^{\prime 2}} d x+\int_{\delta}^{1}\left(\mu Y^{\prime 2}-\frac{Y^{2}}{x^{2}}\right) \zeta^{2} d x+2 \mu \int_{\delta}^{1} Y Y^{\prime} \zeta \zeta^{\prime} d x .
\end{aligned}
$$

But

$$
\begin{equation*}
2 \mu \int_{\delta}^{1} Y Y^{\prime} \zeta \zeta^{\prime} d x=-\mu\left(Y Y^{\prime} \zeta^{2}\right)_{\delta}-\mu \int_{\delta}^{1}\left(Y^{\prime 2}+Y Y^{\prime \prime}\right) \zeta^{2} d x \tag{7.4.5}
\end{equation*}
$$

Combining (7.4.4) and (7.4.5), and observing that $Y$ is a solution of

$$
\mu Y^{\prime \prime}+\frac{Y}{x^{2}}=0
$$

we obtain

$$
\begin{equation*}
J_{\delta}(\eta)=-\mu\left(Y Y^{\prime} \zeta^{2}\right)_{\delta}+\mu \int_{\delta}^{1}\left(Y \zeta^{\prime}\right)^{2} d x \tag{7.4.6}
\end{equation*}
$$

But

$$
Y Y^{\prime} \zeta^{2}=\left(\frac{1}{2}+a\right) \frac{(Y \zeta)^{2}}{x}=\left(\frac{1}{2}+a\right) \frac{\eta^{2}}{x} \rightarrow 0
$$

when $x \rightarrow 0$. Hence, when we make $\delta \rightarrow 0$ in (7.4.6), we obtain

$$
\begin{equation*}
J(\eta)=\mu \int_{0}^{1}\left(Y \zeta^{\prime}\right)^{2} d x \tag{7.4.7}
\end{equation*}
$$

which is positive unless the integrand is nul, i.e. unless $\zeta=0$.
This proves the theorem. The condition $\mu>4$ was required to make $Y^{\prime}$ belong to $L^{2}$. We have however reduced the theorem to dependence on the identity (7.4.7). Since

$$
Y \zeta^{\prime}=\eta^{\prime}-\frac{\frac{1}{2}+a}{x} \eta=y^{\prime}-\frac{\frac{1}{2}+a}{x} y
$$

(7.4.7) and (7.4.3) give

$$
\begin{equation*}
\text { 8) } J(y)=\int_{0}^{1}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x=\frac{2}{1-2 a}+\mu \int_{0}^{1}\left(y^{\prime}-\frac{\frac{1}{2}+a}{x} y\right)^{2} d x \tag{7.4.8}
\end{equation*}
$$

and here $Y$ has disappeared. Since both sides of (7.4.8) are continuous in $\mu$, we may now include the case $\mu=4, a=0$. The identity may be verified directly by partial integration when its form has been discovered (though some care is required about convergence at the lower limit).

We are now in a position to prove Theorem 253. If we write

$$
x=X / \xi, \quad c y(x)=Y(X)
$$

and then replace $X, Y$ again by $x, y$, we obtain

$$
\begin{equation*}
\int_{0}^{\xi}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x=\frac{2}{1-2 a} \frac{c^{2}}{\xi}+\mu \int_{0}^{\xi}\left(y^{\prime}-\frac{\frac{1}{2}+a}{x} y\right)^{2} d x \tag{7.4.9}
\end{equation*}
$$

where now $y(0)=0, y(\xi)=c$. If $y^{\prime}$ is $L^{2}$ in ( $0, \infty$ ), $c=o\left(\xi^{\frac{1}{2}}\right)$ for large $\xi$, a and the first term on the right tends to 0 when $\xi \rightarrow \infty$. Making $\xi \rightarrow \infty$, and supposing $\mu=4$, we obtain

$$
\int_{0}^{\infty}\left(4 y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x=4 \int_{0}^{\infty}\left(y^{\prime}-\frac{y}{2 x}\right)^{2} d x
$$

This formula, which makes Theorem 253 intuitive, is valid whenever $y^{\prime}$ is $L^{2}$, and may of course be verified directly ${ }^{\text {b }}$.
7.5. Second proof of Theorem 254. In our second proof we make explicit the variational theory which underlies the first.

Suppose that $y=Y(x)$, or $E$, is the extremal through the endpoints $P_{0}$ and $P_{1}$, and that

$$
\begin{equation*}
y=y(x, \alpha) \tag{7.5.1}
\end{equation*}
$$

or $E(\alpha)$, is a family of extremals containing $E$ and depending on a parameter $\alpha$. Suppose further either that
(i) $E(\alpha)$ covers up a region surrounding $E$ in a (1, l) manner, so that just one extremal passes through every point of the region, and $\alpha$ is a one-valued function of $x$ and $y$; or that
(ii) every curve of $E(\alpha)$ passes through $P_{0}$, so that $y\left(x_{0}, \alpha\right)$ is independent of $\alpha$, but condition (i) is satisfied in all other respects.
a Theorem 223.
${ }^{\text {b }}$ Grandjot (1) gives a number of somewhat similar identities for series.

In these circumstances ${ }^{\text {a }}$ (7.5.1) is said to define a field of extremals including $E$.

The slope

$$
y^{\prime}(x, \alpha)
$$

of the extremal through a point $P$ of the field may be expressed as a one-valued function

$$
p=p(x, y)
$$

of $x$ and $y$. Hilbert's 'invariant integral' is

$$
J^{*}(C)=\int_{C}\left\{\left(F-p F_{p}\right) d x+F_{p} d y\right\} .
$$

Here $F$ and $F_{p}$ are the values of $F\left(x, y, y^{\prime}\right)$ and $F_{y^{\prime}}\left(x, y, y^{\prime}\right)$ when $y^{\prime}$ is replaced by $p$, and the integral is taken along any curve $C$ which lies in the region covered by the field.

The fundamental properties of Hilbert's integral are as follows.
(i) $J^{*}(C)$ depends only on the ends $Q, R$ of $C$; in other words

$$
\left(F-p F_{p}\right) d x+F_{p} d y=d W
$$

is a perfect differential, and

$$
J^{*}(C)=W_{R}-W_{Q} .
$$

(ii) If $C$ is the extremal $E$, then

$$
J^{*}(E)=\int_{E} F d x=J(E),
$$

say. It follows that, if $C$ runs from $P_{0}$ to $P_{1}$, then

$$
\begin{aligned}
J(C)-J(E)= & J(C)-J^{*}(E)=J(C)-J^{*}(C) \\
= & \int_{C} F\left(x, y, y^{\prime}\right) d x \\
& -\int_{C}\left\{\left(F(x, y, p)-p F_{p}(x, y, p)\right) d x+F_{p}(x, y, p) d y\right\} \\
= & \int_{C} \delta\left(x, y, p, y^{\prime}\right) d x
\end{aligned}
$$

where

$$
\varepsilon\left(x, y, p, y^{\prime}\right)=F\left(x, y, y^{\prime}\right)-F(x, y, p)-\left(y^{\prime}-p\right) F_{p}(x, y, p) .
$$

a With the addition of certain conditions concerning the differentiability of $\alpha(x, y)$ which it is unnecessary to repeat here: see Bolza (1,95-105).

Here $y^{\prime}$ is the slope of $C$ at any point and $p$ the slope of the extremal through the point; and $\varepsilon$ is Weierstrass's 'excessfunction'. If $\varepsilon>0$ whenever $y^{\prime} \neq p$, then

$$
J(E)<J(C)
$$

and $E$ gives a true minimum of $J$.
In the present case we take
as $E(\alpha)$. We find

$$
\begin{gathered}
p=\alpha\left(\frac{1}{2}+a\right) x^{-\frac{1}{2}+a}=\left(\frac{1}{2}+a\right) \frac{y}{x} \\
F=\mu p^{2}-\frac{y^{2}}{x^{2}}, \quad F_{p}=\frac{2 y}{\left(\frac{1}{2}-a\right) x}, \quad F-p F_{p}=-\frac{y^{2}}{\left(\frac{1}{2}-a\right) x^{2}}, \\
J^{*}=\frac{1}{\frac{1}{2}-a} \int\left(-\frac{y^{2}}{x^{2}} d x+\frac{2 y}{x} d y\right)=\int d W
\end{gathered}
$$

where

$$
y=\alpha x^{\frac{1}{2}+a}
$$

Here $\quad \varepsilon=\mu y^{\prime 2}-\mu p^{2}-\left(y^{\prime}-p\right) 2 \mu p=\mu\left(y^{\prime}-p\right)^{2}>0$
unless $y^{\prime}=p$. The identity

$$
J(C)-J(E)=\int_{C} \varepsilon d x
$$

reduces to

$$
\int_{0}^{1}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x-\frac{2}{1-2 a}=\mu \int_{0}^{1}\left\{y^{\prime}-\left(\frac{1}{2}+a\right) \frac{y}{x}\right\}^{2} d x
$$

which is (7.4.8).
This argument shows the genesis of (7.4.8), but does not prove it, for two reasons. In the first place, $F$ has a singularity, and the theory of the field breaks down, for $x=0$. Secondly, the theory presupposes the continuity of $y^{\prime}$.

In order to dispose of the first difficulty, we may take $P_{0}$ and $P_{1}$ to be $\left(\delta, \delta^{\frac{1}{2}+a}\right)$ and $(1,1)$. The theory then gives the identity

$$
\int_{\delta}^{1}\left(\mu y^{\prime 2}-\frac{y^{2}}{x^{2}}\right) d x=\frac{2\left(1-\delta^{2 a}\right)}{1-2 a}+\mu \int_{\delta}^{1}\left\{y^{\prime}-\left(\frac{1}{2}+a\right) \frac{y}{x}\right\}^{2} d x
$$

and we obtain (7.4.8), for continuous $y^{\prime}$, by making $\delta$ tend to zero.

When (7.4.8) is proved for continuous $y^{\prime}$, it may be extended
to general $y^{\prime}$ of $L^{2}$ by standard processes of approximation. We deal with this point in the next section, in a different problem where we have no alternative elementary proof.

Similar considerations lead to the identity

$$
\begin{align*}
& \int_{0}^{1}\left(\mu y^{\prime k}-\frac{y^{k}}{x^{k}}\right) d x=\frac{1}{(k-1)(1-\lambda)}  \tag{7.5.2}\\
& \quad+\mu \int_{0}^{1}\left\{y^{k}-\left(\frac{\lambda y}{x}\right)^{k}-k\left(y^{\prime}-\frac{\lambda y}{x}\right)\left(\frac{\lambda y}{x}\right)^{k-1}\right\} d x .^{\mathrm{a}}
\end{align*}
$$

Here

$$
y(0)=0, \quad y(1)=1, \quad y^{\prime} \geqq 0, \quad k>1, \quad \mu>\left(\frac{k}{k-1}\right)^{k}=K
$$

and $\lambda$ is the (unique) root of

$$
\begin{equation*}
\mu(k-1) \lambda^{k-1}(\lambda-1)+1=0 \tag{7.5.3}
\end{equation*}
$$

which lies between $1 / k^{\prime}$ and 1 . When the form of (7.5.2) has been determined we may put

$$
\mu=K
$$

where $\lambda=1 / k^{\prime}$ and $\mu \lambda^{k}=1$. We thus obtain

$$
\begin{aligned}
\int_{0}^{1}\left(K y^{\prime k}-\right. & \left.-\frac{y^{k}}{x^{k}}\right) d x \\
& =\frac{k}{k-1}+K \int_{0}^{1}\left\{y^{\prime k}-\left(\frac{\lambda y}{x}\right)^{k}-k\left(y^{\prime}-\frac{\lambda y}{x}\right)\left(\frac{\lambda y}{x}\right)^{k-1}\right\} d x
\end{aligned}
$$

It may be verified directly, by partial integration, that this is true whenever $y^{\prime}$ is $L^{k}$; and we can prove as in $\S 7.4$ that the identity remains true when the upper limit 1 is replaced by $\infty$ and the term $k /(k-1)$ is omitted. Since, by Theorem 41,

$$
a^{k}-b^{k}>k(a-b) b^{k-1}
$$

for all positive $a, b$, we thus obtain a proof of a theorem (Theorem 327) which will be stated explicitly, and proved in an entirely different manner, in §9.8.

Incidentally we obtain
255. If

$$
k>1, \quad \mu>\left(\frac{k}{k-1}\right)^{k}=K
$$

${ }^{\text {a }}$ The theory of the field gives the form of the identity, which may then be verified independently. The limitation to curves for which $y^{\prime} \geqq 0$ would introduce another slight complication into a properly variational proof.
$y(0)=0, y(1)=1$, and $y^{\prime}$ is $L^{k}$, then

$$
J=\int_{0}^{1}\left(\mu y^{\prime k}-\frac{y^{k}}{x^{k}}\right) d x \geqq \frac{1}{(k-1)(1-\lambda)}
$$

where $\lambda$ is the root of $(7.5 .3)$ between $1 / k^{\prime}$ and 1 .
7.6. Further examples illustrative of variational methods. It is difficult to distinguish at all precisely between 'elementary' and 'variational' proofs, since there are many proofs of intermediate types. We give a selection of such proofs, worked out with varying degrees of detail, in this and the succeeding sections.
(I). 256. If $y(0)=0$ and $2 k$ is an even positive integer, then

$$
\begin{equation*}
\int_{0}^{1} y^{2 k} d x \leqq C \int_{0}^{1} y^{\prime 2 k} d x \tag{7.6.1}
\end{equation*}
$$

where

$$
\begin{equation*}
C=\frac{1}{2 k-1}\left(\frac{2 k}{\pi} \sin \frac{\pi}{2 k}\right)^{2 k} \tag{7.6.2}
\end{equation*}
$$

There is equality only for a certain hyperelliptic curve.
(i) We suppose first that $y(1) \neq 0$, in which case we may take $y(1)=1$, and consider

$$
J(y)=\int_{0}^{1}\left(C y^{\prime 2 k}-y^{2 k}\right) d x
$$

Euler's equation is

$$
(2 k-1) C y^{\prime 2 k-2} y^{\prime \prime}+y^{2 k-1}=0
$$

which gives

$$
(2 k-1) C y^{\prime 2 k}=C^{\prime}-y^{2 k}
$$

where $C^{\prime}$ is a constant of integration.
There is one extremal which passes through $(0,0)$ and $(1,1)$ and cuts $x=1$ at right angles. In fact, if we take $C^{\prime}=1$, then $y^{\prime}$ vanishes when $y=1$. Also

$$
x=\{(2 k-1) C\}^{1 / 2 k} \int \frac{d y}{\left(1-y^{2 k}\right)^{1 / 2 k}}
$$

and, since

$$
\int_{0}^{1} \frac{d y}{\left(1-y^{2 k}\right)^{1 / 2 k}}=\frac{1}{2 k} \int_{0}^{1}(1-u)^{-1 / 2 k} u^{1 / 2 k-1} d u=\frac{\pi}{2 k} \operatorname{cosec} \frac{\pi}{2 k},
$$

there is an extremal of this type which passes through $(0,0)$ and $(1,1)$. If we denote this extremal by $Y$, then
$J(Y)=\left[C Y Y^{\prime 2 k-1}\right]_{0}^{1}-\int_{0}^{1} Y\left\{(2 k-1) C Y^{\prime 2 k-2} Y^{\prime \prime}+Y^{2 k-1}\right\} d x=0$, since $Y^{\prime}(1)=0$.

To prove the theorem, we must show that $y=Y$ gives a strong minimum; and this follows easily from the general theory. The extremal is a rising curve of the same general form as the curve $y=\sin \frac{1}{2} \pi x$, to which it reduces when $k=1$. The curve $y=\alpha Y$ is also an extremal; the family $y=\alpha Y$ defines a field in the sense of §7.5; and the excess-function

$$
\varepsilon=y^{\prime 2 k}-p^{2 k}-2 k\left(y^{\prime}-p\right) p^{2 k-1}
$$

is positive. Hence the standard conditions for a minimum are satisfied. This proof is genuinely 'variational', and (in view of the trouble of calculating the slope-function $p$ explicitly) it might be difficult to find a more elementary proof.

There is however one point in the proof which demands an additional remark. The 'general theory' assumes that $y$ ' is continuous, and it may not be obvious how its conclusions, in particular in regard to the uniqueness of the solution, are extended to the more general $y$ considered here.

Let us denote by $I^{2 p}$ the class of integrals $y$ of functions of $L^{2 p}$, by $I^{*}$ the class of integrals $y^{*}$ of continuous functions. The general theory shows that

$$
\begin{equation*}
J\left(y^{*}\right)>J(Y) \tag{7.6.3}
\end{equation*}
$$

for a $y^{*}$ different from $Y$, while we require the same result for any $y$ of $I^{2 p}$. We can approximate to a $y$ of $I^{2 p}$, different from $Y$, by a sequence of functions $y^{*}$, in such a manner that

$$
J(y)=\lim J\left(y^{*}\right)
$$

but all that then follows from (7.6.3) is

$$
J(y) \geqq J(Y),
$$

the strict sign of inequality being lost in the passage to the limit.
The difficulty disappears if we look at the question differently. The general theory proves not only the inequality (7.6.3) but also the identity

$$
J\left(y^{*}\right)-J(Y)=\int \varepsilon\left(y^{*}\right) d x
$$

where $હ\left(y^{*}\right)$ is the excess-function corresponding to $y^{*}$ and the field $y=\alpha Y$. Approximating to $y$ by an appropriate sequence of $y^{*}$, we replace this identity by

$$
J(y)-J(Y)=\int \varepsilon(y) d x
$$

and the integral is positive unless $y=Y$.
(ii) The case in which $y(1)=0$ may be discussed similarly, since $y=0$ is then an extremal satisfying the conditions, and is included in the field $y=\alpha Y$ used in (i).

The proof might have been arranged differently if we had made no hypothesis about the value of $y(1)$. The problem is then one with a 'variable end-point', that of minimising $J(y)$ for curves drawn from the origin to meet the line $x=1$. The extremals cut this line 'transversally' (in this case orthogonally), and all the curves $y=\alpha Y$ satisfy this condition. The general theory shows that all the extremals give the same value of $J(y)$, and this value must be 0 , since it is 0 when $\alpha=0$. ${ }^{\text {a }}$
7.7. Further examples: Wirtinger's inequality. (II) Let us consider more particularly the case $k=1$ of (I). Changing the limits, the result is that

$$
\begin{equation*}
\int_{0}^{\frac{1}{2} \pi} y^{2} d x<\int_{0}^{\frac{1}{2} \pi} y^{\prime 2} d x \tag{7.7.1}
\end{equation*}
$$

if $y(0)=0$ and $y$ is not a multiple of $\sin x$.
The general theory suggests that there is an identity of the type
or

$$
\begin{gathered}
\int_{0}^{\frac{1}{2} \pi}\left(y^{\prime 2}-y^{2}\right) d x=\int_{0}^{\frac{1}{2} \pi}\left\{y^{\prime}-y \psi(x)\right\}^{2} d x \\
\int_{0}^{\frac{1}{2} \pi}\left\{y^{2}\left(1+\psi^{2}\right)-2 y y^{\prime} \psi\right\} d x=0
\end{gathered}
$$

This will plainly be true if

$$
y^{2}\left(1+\psi^{2}\right) d x-2 y \psi d y
$$

is an exact differential $d z$, and $z$ vanishes at the limits; and this requires

$$
-\psi^{\prime}=1+\psi^{2}, \quad \psi=-\tan (x+k)
$$

a We owe these remarks to Prof. Bliss.
in which case $z=-\psi y^{2}$. If we take $k=\frac{1}{2} \pi, \psi=\cot x$, then $z$ vanishes when $x=\frac{1}{2} \pi$; and since $y^{\prime}$ is $L^{2}$ and so ${ }^{\text {a }} y=o\left(x^{\frac{1}{2}}\right), z$ also vanishes when $x=0$. We thus obtain

$$
\begin{equation*}
\int_{0}^{\frac{1}{2} \pi}\left(y^{\prime 2}-y^{2}\right) d x=\int_{0}^{\frac{1}{2} \pi}\left(y^{\prime}-y \cot x\right)^{2} d x \tag{7.7.2}
\end{equation*}
$$

which makes (7.7.1) intuitive.
A slight modification of (7.7.2) leads to
257. If $y(0)=y(\pi)=0$ and $y^{\prime}$ is $L^{2}$, then

$$
\int_{0}^{\pi} y^{2} d x<\int_{0}^{\pi} y^{\prime 2} d x
$$

unless

$$
y=C \sin x
$$

For $y=o\left(x^{\frac{1}{2}}\right)$ for small $x$, and $y=o\left\{(\pi-x)^{\frac{1}{2}}\right\}$ for $x$ near $\pi$, so that $y^{2} \cot x$ vanishes at both limits. Hence

$$
\begin{equation*}
\int_{0}^{\pi}\left(y^{\prime 2}-y^{2}\right) d x=\int_{0}^{\pi}\left(y^{\prime}-y \cot x\right)^{2} d x \tag{7.7.3}
\end{equation*}
$$

Another modification of (7.7.2) leads to a more interesting theorem due to Wirtinger ${ }^{\text {b }}$.
258. If $y$ has the period $2 \pi, y^{\prime}$ is $L^{2}$, and

$$
\begin{gather*}
\int_{0}^{2 \pi} y d x=0  \tag{7.7.4}\\
\int_{0}^{2 \pi} y^{2} d x<\int_{0}^{2 \pi} y^{\prime 2} d x \\
y=A \cos x+B \sin x
\end{gather*}
$$

then
unless
We cannot write down at once an identity similar to (7.7.2) or (7.7.3), but with $0,2 \pi$ as limits, because $y \cot x$ will usually have infinities in the range of integration. We may however argue as follows ${ }^{c}$.
a Theorem 222
b See Blaschke (1, 105). The most immediate proof is by an application of Parseval's Theorem to the Fourier developments

$$
y \sim \frac{1}{2} a_{0}+\Sigma\left(a_{n} \cos n x+b_{n} \sin n x\right), \quad y^{\prime} \sim \Sigma\left(n b_{n} \cos n x-n a_{n} \sin n x\right)
$$

(with $a_{0}=0$ ).
c The proof which follows was communicated to us by Dr Hans Lewy.

The function

$$
z(x)=y(x+\pi)-y(x)
$$

has opposite signs for the arguments $x$ and $x+\pi$, and therefore vanishes at least once in $(0, \pi)$. We suppose that $z(\alpha)=0$, where $0 \leqq \alpha<\pi$, and write $y(\alpha)=a$. Since $y^{\prime}$ is $L^{2},(y-a)^{2} \cot (x-\alpha)$ vanishes for $x=\alpha$ and $x=\alpha+\pi ;$ a and

$$
\begin{aligned}
\int_{0}^{2 \pi}\left[y^{\prime 2}-(y-a)^{2}-\left\{y^{\prime}-(y-a)\right.\right. & \left.\cot (x-\alpha)\}^{2}\right] d x \\
& =\left[(y-a)^{2} \cot (x-\alpha)\right]_{0}^{2 \pi}=0
\end{aligned}
$$

Hence, using (7.7.4), we obtain

$$
\int_{0}^{2 \pi}\left(y^{\prime 2}-y^{2}\right) d x=2 \pi a^{2}+\int_{0}^{2 \pi}\left\{y^{\prime}-(y-a) \cot (x-\alpha)\right\}^{2} d x
$$

which is positive unless $a=0$ and

$$
y^{\prime} \equiv y \cot (x-\alpha), \quad y=C \sin (x-\alpha)
$$

There is a special interest in Theorem 258 because the proof of the classical isoperimetric property of the circle may be based upon it. We consider a simple closed curve $C$ whose area is $A$ and whose perimeter is $L$, and take

$$
\phi=\frac{2 \pi s}{L}
$$

where $s$ is the arc of the curve, as parameter, so that

$$
x=x(\phi), \quad y=y(\phi) \quad(0 \leqq \phi \leqq 2 \pi)
$$

We suppose for simplicity that $x^{\prime}$ and $y^{\prime}$ are continuous; the proof is valid for more general $x, y$. We may also suppose without loss of generality that the centre of gravity of the perimeter lies on the axis of $x$, so that

$$
\int_{0}^{2 \pi} y d \phi=0
$$

We have then

$$
\left(\frac{d x}{d \phi}\right)^{2}+\left(\frac{d y}{d \phi}\right)^{2}=\left\{\left(\frac{d x}{d s}\right)^{2}+\left(\frac{d y}{d s}\right)^{2}\right\} \frac{L^{2}}{4 \pi^{2}}=\frac{L^{2}}{4 \pi^{2}}
$$

a Using Theorem 222 as in § 7.4 and above.

$$
\begin{aligned}
\frac{L^{2}}{2 \pi}-2 A & =\int_{0}^{2 \pi}\left\{\left(\frac{d x}{d \bar{\phi}}\right)^{2}+\left(\frac{d y}{d \phi}\right)^{2}\right\} d \phi+2 \int_{0}^{2 \pi} y \frac{d x}{d \bar{\phi}} d \phi^{\mathrm{a}} \\
& =\int_{0}^{2 \pi}\left(\frac{d x}{d \phi}+y\right)^{2} d \phi+\int_{0}^{2 \pi}\left\{\left(\frac{d y}{d \phi}\right)^{2}-y^{2}\right\} d \phi \\
& \geqq \int_{0}^{2 \pi}\left\{\left(\frac{d y}{d \phi}\right)^{2}-y^{2}\right\} d \phi \geqq 0,
\end{aligned}
$$

by Theorem 258. There will be inequality unless

$$
y=A \cos \phi+B \sin \phi
$$

and also $\quad x=-\int y d \phi=-A \sin \phi+B \cos \phi+C$,
when the curve is a circle ${ }^{b}$.
7.8. An example involving second derivatives. (III) It is knownc that, if $f$ has a second derivative for $x \geqq 0$, and $\mu_{0}, \mu_{1}, \mu_{2}$ are the upper bounds of $|f|,\left|f^{\prime}\right|,\left|f^{\prime \prime}\right|$, then

$$
\mu_{1}{ }^{2} \leqq 4 \mu_{0} \mu_{2}
$$

This suggests that there may be a corresponding relation between the integrals

$$
J_{0}=\int_{0}^{\infty}|f|^{p} d x, \quad J_{1}=\int_{0}^{\infty}\left|f^{\prime}\right|^{p} d x, \quad J_{2}=\int_{0}^{\infty}\left|f^{\prime \prime}\right|^{p} d x
$$

where $p \geqq 1$. The next theorem settles this question in the case $p=2$.
259. If $y$ and $y^{\prime \prime}$ are $L^{2}$ in $(0, \infty)^{\text {d }}$, then

$$
\left(\int_{0}^{\infty} y^{\prime 2} d x\right)^{2}<4 \int_{0}^{\infty} y^{2} d x \int_{0}^{\infty} y^{\prime \prime 2} d x
$$

unless $y=A Y(B x)$, where

$$
Y=e^{-\frac{1}{2} x} \sin (x \sin \gamma-\gamma) \quad\left(\gamma=\frac{1}{3} \pi\right),
$$

when there is equality.
a Or $-2 \int_{0}^{2 \pi} y \frac{d x}{d \phi} d \phi$, according to the sense of the variation of $s$ as we pass in a given sense round the curve.
b The proof is in principle that of Hurwitz (2), but differs (a) in that we do not use the theory of Fourier series and (b) in our unsymmetrical treatment of $x$ and $y$.
c Landau (2, 3).
d In accordance with the convention of $\S 7.1, y^{\prime}$ is the integral of $y^{\prime \prime}$ and $y$ of $y^{\prime}$.

If we consider the 'isoperimetrical' problem of the Calculus of Variations defined by
${ }^{'} J_{1}=\int_{0}^{\infty} y^{\prime 2} d x$ maximum, $J_{0}=\int_{0}^{\infty} y^{2} d x$ and $J_{2}=\int_{0}^{\infty} y^{\prime \prime 2} d x$ given', we have to form Euler's equation for

$$
\int_{0}^{\infty}\left(y^{\prime 2}-\lambda y^{2}-\mu y^{\prime \prime 2}\right) d x
$$

It is a linear equation of the form

$$
a y^{\prime \prime \prime \prime}+b y^{\prime \prime}+c y=0
$$

whose solutions are linear combinations of real or complex exponentials. When we try to choose the parameters in the most advantageous way, we are led to consider the function $Y$.

It seems difficult to complete the proof on these lines by use of the general theory. We shall deduce Theorem 259 from the simpler theorem which follows.
260. Under the conditions of Theorem 259,

$$
J(y)=\int_{0}^{\infty}\left(y^{2}-y^{\prime 2}+y^{\prime \prime 2}\right) d x>0
$$

unless $y=A Y$, when there is equality.
We give several proofs of this theorem to illustrate differences of method. The first two are, as they stand, elementary; the third, of which we give only an outline, makes explicit the variational theory which lies behind the other two. We begin by an observation which is necessary in any case, viz. that $J_{1}$ is finite.

To prove this, we have

$$
\begin{equation*}
\int_{0}^{x} y^{\prime 2} d x=\left[y y^{\prime}\right]_{0}^{X}-\int_{0}^{X} y y^{\prime \prime} d x \tag{7.8.1}
\end{equation*}
$$

Since $J_{0}$ and $J_{2}$ are finite, the last integral tends to a finite limita when $X \rightarrow \infty$. If $J_{1}$ were infinite, $y y^{\prime}$ and $a$ fortiori

$$
y^{2}=2 \int y y^{\prime} d x
$$

would tend to infinity, which is impossible on account of the convergence of $J_{0}$. Hence $J_{1}$ is finite, and all three terms in (7.8.1)

[^46]tend to limits. In particular $y y^{\prime}$ tends to a limit, which can only be 0 (again on account of the convergence of $J_{0}$ ).
(1) Our first proof proceeds rather on the lines of §7.4. It is easily verified that
$$
Y+Y^{\prime}+Y^{\prime \prime}=0, \quad Y(0)+Y^{\prime}(0)=0, \quad Y^{\prime \prime}(0)=0
$$
and
$$
J(Y)=0 .
$$

We now write

$$
y=z+c Y
$$

choosing $c$ so that $z(0)=0$. Then $z, z^{\prime}$, and $z^{\prime \prime}$ are $L^{2}, z$ and $z^{\prime}$ are $o\left(x^{\frac{1}{2}}\right)$ for large $x,{ }^{\mathrm{b}}$ and $z z^{\prime} \rightarrow 0$.

Now

$$
J(y)=J(z)+2 c K(Y, z)+c^{2} J(Y)
$$

where

$$
\begin{aligned}
K & =\int_{0}^{\infty}\left(Y z-Y^{\prime} z^{\prime}+Y^{\prime \prime} z^{\prime \prime}\right) d x \\
& =-\int_{0}^{\infty}\left(Y^{\prime}+Y^{\prime \prime}\right) z d x-\int_{0}^{\infty}\left(Y+Y^{\prime}\right) z^{\prime \prime} d x-\int_{0}^{\infty} Y^{\prime} z^{\prime} d x \\
& =\int_{0}^{\infty}\left(Y+Y^{\prime}+Y^{\prime \prime}\right) z^{\prime} d x=0
\end{aligned}
$$

Hence $J(y)=J(z)$, and it is enough to prove that $J(z)>0$ unless $z=0$. But, since $z(0)=0$ and $z z^{\prime} \rightarrow 0$ when $x \rightarrow \infty$,

$$
\int_{0}^{\infty} z^{\prime 2} d x=-\int_{0}^{\infty} z z^{\prime \prime} d x
$$

and so

$$
J(z)=\int_{0}^{\infty}\left(z^{2}+z z^{\prime \prime}+z^{\prime \prime 2}\right) d x>0
$$

unless $z=0$. This proves Theorem 260.
(2) We may try (following the lines of § 7.5 and $\S 7.7$ ) to reduce Theorem 260 to dependence upon an identity. For this, we make

$$
\begin{equation*}
\left\{y^{2}-y^{\prime 2}+y^{\prime \prime 2}-\left(y^{\prime \prime}+\phi y+\psi y^{\prime}\right)^{2}\right\} d x \tag{7.8.2}
\end{equation*}
$$

an exact differential; and the simplest choice of $\phi$ and $\psi$ is $\phi=\psi=1$, when (7.8.2) reduces to

$$
-d\left(y+y^{\prime}\right)^{2}
$$

Thus

$$
\int_{0}^{x}\left\{y^{2}-y^{\prime 2}+y^{\prime \prime 2}-\left(y+y^{\prime}+y^{\prime \prime}\right)^{2}\right\} d x=-\left[\left(y+y^{\prime}\right)^{2}\right]_{0}^{x}
$$

Since $J_{0}, J_{1}$, and $J_{2}$ are finite, the left-hand side has a limit when $X \rightarrow \infty .^{\text {a }}$ Also $y y^{\prime} \rightarrow 0$, and so $y^{2}+y^{\prime 2}$ tends to a limit, which can only be 0 . It follows that

$$
\begin{equation*}
\int_{0}^{\infty}\left(y^{2}-y^{\prime 2}+y^{\prime \prime 2}\right) d x=\left\{y(0)+y^{\prime}(0)\right\}^{2}+\int_{0}^{\infty}\left(y+y^{\prime}+y^{\prime \prime}\right)^{2} d x \tag{7.8.3}
\end{equation*}
$$

which is positive unless

$$
\begin{gathered}
y^{\prime \prime}+y^{\prime}+y=0 \\
y^{\prime}(0)+y(0)=0 .
\end{gathered}
$$

and
These conditions give $y=A Y$.
(3) The underlying variational theory is a little more complex than that of §7.5. If we put

$$
\begin{gather*}
y^{\prime}=z  \tag{7.8.4}\\
J(y)=J(y, z)=\int_{0}^{\infty}\left(y^{2}-y^{\prime 2}+z^{\prime 2}\right) d x \tag{7.8.5}
\end{gather*}
$$

and suppose that

$$
\begin{equation*}
y(0)=1, \quad z(0)=\zeta, \quad y(\infty)=0, \quad z(\infty)=0, \tag{7.8.6}
\end{equation*}
$$

then the problem is a 'Lagrange problem', viz. that of minimising $J(y, z)$ when subject to (7.8.4) and (7.8.6).

The extremals of the field, in space $(x, y, z)$ are given ${ }^{\text {b }}$ by

$$
\begin{equation*}
\frac{\partial \Phi}{\partial y}-\frac{d}{d x}\left(\frac{\partial \Phi}{\partial y^{\prime}}\right)=0, \quad \frac{\partial \Phi}{\partial z}-\frac{d}{d x}\left(\frac{\partial \Phi}{\partial z^{\prime}}\right)=0, \tag{7.8.7}
\end{equation*}
$$

where

$$
\Phi=F-\lambda\left(y^{\prime}-z\right)=y^{2}-y^{\prime 2}+z^{\prime 2}-\lambda\left(y^{\prime}-z\right),
$$

and $\lambda$ is a function of $x$ defined by the equations themselves. In this case the equations (7.8.7) reduce to

$$
\begin{equation*}
2 y+\frac{d}{d x}\left(2 y^{\prime}+\lambda\right)=0, \quad \lambda=\frac{d}{d x}\left(2 z^{\prime}\right), \tag{7.8.8}
\end{equation*}
$$

and from these and (7.8.4) we find

$$
\begin{equation*}
y^{\prime \prime \prime \prime}+y^{\prime \prime}+y=0 \tag{7.8.9}
\end{equation*}
$$

a The integrals

$$
\int_{0}^{\infty} y y^{\prime} d x, \quad \int_{0}^{\infty} y^{\prime} y^{\prime \prime} d x, \quad \int_{0}^{\infty} y y^{\prime \prime} d x
$$

being convergent, by Theorem 181.
b For an account of the general theory see Bliss (2).

The most general solution of (7.8.9), for which $y$ and $z$ vanish at infinity, is

$$
\begin{equation*}
y=a e^{-\rho x}+\bar{a} e^{-\bar{\rho} x}, \tag{7.8.10}
\end{equation*}
$$

where $\rho=e^{\frac{\xi}{\pi r i}}$ and the bar denotes the conjugate.
The equations (7.8.10) and

$$
\begin{gather*}
z=-a \rho e^{-\rho x}-\bar{a} \bar{\rho} e^{-\bar{\rho} x}  \tag{7.8.11}\\
\lambda=2 a e^{-\rho x}+2 \bar{a} e^{-\bar{\rho} x} \tag{7.8.12}
\end{gather*}
$$

define a two-parameter 'field' of extremals through ( $\infty, 0,0$ ). The 'slope functions' $p, q$, and the 'multiplier' $\lambda$ of the field are the functions

$$
p(x, y, z), \quad q(x, y, z), \quad \lambda(x, y, z),
$$

obtained by expressing the slopes of the extremal, and the function $\lambda$ defined by (7.8.8), in terms of $x, y, z$. A straightforward calculation gives

$$
\begin{gather*}
p=z, \quad q=-y-z,  \tag{7.8.13}\\
\lambda=2 y .{ }^{a} \tag{7.8.14}
\end{gather*}
$$

The analogue of the Hilbert integral is

$$
J^{*}=\int\left[\left(\Phi-p \Phi_{p}-q \Phi_{q}\right) d x+\Phi_{p} d y+\Phi_{q} d z\right] .
$$

This integral has properties corresponding to those of $\S 7.5$; it is independent of the path between its end-points, and its value along an extremal is the same as that of

$$
J=\int F d x
$$

Also, if $E$ is the extremal, and $C$ any other curve, joining the end-points, we have

$$
\begin{equation*}
J_{C}-J_{E}=J_{C}-J_{E}^{*}=J_{C}-J_{C}^{*}=\int_{C} \varepsilon d x \tag{7.8.15}
\end{equation*}
$$

where $\varepsilon$ is the excess-function, defined here by

$$
\varepsilon=\Phi\left(x, y, z, y^{\prime}, z^{\prime}, \lambda\right)-\Phi(x, y, z, p, q, \lambda)-\left(y^{\prime}-p\right) \Phi_{p}-\left(z^{\prime}-q\right) \Phi_{q} .
$$

${ }^{\text {a }} p, q, \lambda$ are in the first instance functions of $x$ and the parameters $a, \bar{a}$ of the extremal. Here, in particular,

$$
\begin{gathered}
y^{\prime}=-a \rho e^{-\rho x}-\bar{a} \bar{\rho} e^{-\bar{\rho} x}=z, \quad z^{\prime}=a \rho^{2} e^{-\rho x}+\bar{a} \bar{\rho}^{2} e^{-\bar{\rho} x}=-y-z, \\
\lambda=2 z^{\prime \prime}=-2 a \rho^{3} e^{-\rho x}-2 \bar{a} \bar{\rho}^{3} e^{-\bar{\rho} x}=2 y .
\end{gathered}
$$

In this case it will be found that

$$
\begin{gather*}
\varepsilon=\left(y+z+z^{\prime}\right)^{2}-\left(y^{\prime}-z\right)^{2}  \tag{7.8.16}\\
J_{E}=(1+\zeta)^{2}
\end{gather*}
$$

and
Since $\varepsilon$ reduces to $\left(y+y^{\prime}+z\right)^{2}$ when $y^{\prime}=z$, we obtain

$$
\int_{0}^{\infty}\left(y^{2}-y^{\prime 2}+z^{\prime 2}\right) d x=(1+\zeta)^{2}+\int_{0}^{\infty}\left(y+y^{\prime}+z^{\prime}\right)^{2} d x
$$

which is (7.8.3). We have thus proved Theorem 260 in three different ways.

We have supposed that $y(0) \neq 0$, so that we can take $y(0)=1$. The case in which $y(0)=0$ may be discussed similarly ${ }^{\text {a }}$.

Our object has been only to illustrate a method and not to discuss a difficult general theory, and we have therefore presented the argument very shortly. The following remarks ${ }^{\text {b }}$ may help to make the method intelligible.
(a) The integral $J^{*}$ constructed from a two-parameter set of extremals is not necessarily invariant. Here we can verify the invariance of $J^{*}$ directly; in fact

$$
J^{*}=-\int d(y+z)^{2}
$$

This invariance ensures that our extremals form a 'field'. The 'reason' is to be found in the fact that they pass through a fixed point, $y, z$, and $\lambda$ all vanishing for $x=\infty$.
(b) In this case $\Phi$ is quadratic, and

$$
\varepsilon=\frac{1}{2}\left\{\left(y^{\prime}-p\right)^{2} \Phi_{p p}+2\left(y^{\prime}-p\right)\left(z^{\prime}-q\right) \Phi_{p q}+\left(z^{\prime}-q\right)^{2} \Phi_{q Q}\right\},
$$

which leads immediately to the formula (7.8.16).
(c) Suppose that $E$ and $C$ are defined as above, that $E_{0}$ is the positive axis of $x$, and $L$ an arbitrary curve joining $(0,1, \zeta)$ to the origin. Then

$$
\begin{aligned}
J_{C} & =J_{C}-J_{E_{0}}=J_{C}-J_{E_{0}}^{\star}=J_{L}^{\star}+J_{C}-J_{L+E_{0}}^{\star} \\
& =J_{L}^{\star}+J_{C}-J_{C}^{\star}=J_{L}^{\star}+\int_{0}^{\infty} \varepsilon d x \\
& =(1+\zeta)^{2}+\int_{0}^{\infty}\left(y+y^{\prime}+y^{\prime \prime}\right)^{2} d x .
\end{aligned}
$$

This is (7.8.3); the argument avoids a direct calculation of $J_{E}$. Alternatively we may argue

$$
\begin{gathered}
J_{C}-J_{E}=\int_{0}^{\infty} \varepsilon d x, \\
J_{E}=J_{E}^{\star}=J_{L}^{\star}+J_{E_{0}}^{\star}=(1+\zeta)^{2} .
\end{gathered}
$$

[^47]In order to deduce Theorem 259 from Theorem 260, we apply the last theorem to $y(x / \rho)$ instead of to $y(x)$. We thus obtain

$$
\rho^{4} J_{0}-\rho^{2} J_{1}+J_{2}=\int_{0}^{\infty}\left(\rho^{4} y^{2}-\rho^{2} y^{\prime 2}+y^{\prime \prime 2}\right) d x>0 .
$$

Since this is true for all positive $\rho$, and in particular when $\rho^{2}=J_{1} / 2 J_{0}$, it follows that $J_{1}{ }^{2}<4 J_{0} J_{2}$ unless $y(x / B)=A Y(x)$.
7.9. A simpler problem. It is interesting to observe that the corresponding theorem for the interval $(-\infty, \infty)$ is much more elementary and of a quite different character.
261. If $y$ and $y^{\prime \prime}$ are $L^{2}$ in $(-\infty, \infty)$, then

$$
\left(\int_{-\infty}^{\infty} y^{\prime 2} d x\right)^{2}<\int_{-\infty}^{\infty} y^{2} d x \int_{-\infty}^{\infty} y^{\prime \prime 2} d x
$$

unless $y=0$. The (unit) constant is the best possible.
In fact (as in the proof of Theorem 260) $y y^{\prime} \rightarrow 0$ and

$$
J_{1}{ }^{2}=\left(\int_{-\infty}^{\infty} y^{\prime 2} d x\right)^{2}=\left(-\int_{-\infty}^{\infty} y y^{\prime \prime} d x\right)^{2}<\int_{-\infty}^{\infty} y^{2} d x \int_{-\infty}^{\infty} y^{\prime \prime 2} d x=J_{0} J_{2}
$$

To prove the constant best possible, we take $y=\sin x$ for $|x| \leqq n \pi, y=0$ for $|x|>n \pi$, and round off the angles at $x= \pm n \pi$ so as to make $y^{\prime \prime}$ continuous. We can plainly do this with changes in the three integrals which are bounded when $n \rightarrow \infty$, and then each of them differs boundedly from $n \pi$, so that
if $n$ is sufficiently large.

$$
J_{1}^{2}>(1-\epsilon) J_{0} J_{2}
$$

## MISCELLANEOUS THEOREMS AND EXAMPLES

262. If $y(0)=y(1)=0$ and $y^{\prime}$ is $L^{2}$, then

$$
\int_{0}^{1} \frac{y^{2}}{x(1-x)} d x<\frac{1}{2} \int_{0}^{1} y^{\prime 2} d x
$$

unless $y=c x(1-x)$.
[If

$$
J(y)=\int_{0}^{1}\left(\frac{1}{2} y^{\prime 2}-\frac{y^{2}}{x(1-x)}\right) d x
$$

then $(E)$ is

$$
x(1-x) y^{\prime \prime}+2 y=0,
$$

and $y=\alpha x(1-x)$ is an extremal satisfying the conditions whatever be $\alpha$. By varying $\alpha$ we can define a field ${ }^{\text {a round any particular extremal. It will }}$ be found that in this case $J(Y)=0$ and

$$
p=\frac{1-2 x}{x(1-x)} y, \quad W=\frac{1-2 x}{2 x(1-x)} y^{2}, \quad \varepsilon=\frac{1}{2}\left\{y^{\prime}-\frac{1-2 x}{x(1-x)} y\right\}^{2} .
$$

[^48]The underlying identity is

$$
\int_{0}^{1}\left(\frac{1}{2} y^{\prime 2}-\frac{y^{2}}{x(1-x)}\right) d x=\frac{1}{2} \int_{0}^{1}\left(y^{\prime}-\frac{1-2 x}{x(1-x)} y\right)^{2} d x .
$$

The theorem should be compared with Theorem 225. For fuller details, and proofs of both theorems by means of Legendre functions, see Hardy and Littlewood (10).]
263. If $J=\int_{0}^{\infty} y^{2} d x, K=\int_{0}^{\infty} y^{\prime 2} d x$, then
unless $y=a e^{-b x}$.

$$
4 J K>\{y(0)\}^{4}
$$

264. If
(a)

$$
y(-1)=-1, \quad y(1)=1, \quad y^{\prime}(-1)=y^{\prime}(1)=0
$$

and $k$ is a positive integer, then

$$
\int_{-1}^{1}\left(y^{\prime \prime}\right)^{2 k} d x \geqq 2\left(\frac{4 k-1}{2 k-1}\right)^{2 k-1}
$$

with inequality unless

$$
y=\frac{4 k-1}{2 k} x-\frac{2 k-1}{2 k} x^{(4 k-1) /(2 k-1)} .
$$

[This is an example of the theory of $\S 7.8$, simpler than that in the text. In this case

$$
\left.\varepsilon=z^{\prime 2 k}-q^{2 k}-2 k\left(z^{\prime}-q\right) q^{2 k-1} .\right]
$$

265. If $y$ satisfies (a), and has a second derivative $y^{\prime \prime}$ for every $x$ of ( $-1,1$ ), then $\left|y^{\prime \prime}(x)\right|>2$ for some $x$.
[This theorem, which is easily proved directly, corresponds to the limiting case $k=\infty$ of Theorem 264. The extremal curve of Theorem 264 reduces to

$$
y=2 x-x^{2} \operatorname{sgn} x
$$

For this curve $y^{\prime}=2(1-|x|)$, and $y^{\prime \prime}=-2 \operatorname{sgn} x$ except for $x=0$. There is no second derivative at the origin.]
266. If $y$ is $L^{2}, z^{\prime}=y$, and
then

$$
y(0)=y(2 \pi)=z(0)=z(2 \pi)=0,
$$

$$
\int_{0}^{2 \pi}\left(y^{\prime 2}-y^{2}\right) d x=\int_{0}^{2 \pi}\left\{y^{\prime}+\frac{(x \cos x-\sin x) y+(1-\cos x) z}{2-2 \cos x-x \sin x}\right\}^{2} d x
$$

[This identity, which gives another (though less simple) proof of Theorem 258, is the result of treating Wirtinger's inequality as a case of Lagrange's problem, on the lines of § 7.8.]
267.

$$
\begin{gathered}
\int_{0}^{\infty}\left(y^{2}+2 y^{\prime 2}+y^{\prime \prime 2}\right) d x>\frac{3}{2}\{y(0)\}^{2} \\
y=C e^{-x}(x+2)
\end{gathered}
$$

268. If $k \geqq 1$ and $y$ and $y^{\prime \prime}$ are $L^{k}$ in $(-\infty, \infty)$ or $(0, \infty)$, then

$$
\left(\int\left|y^{\prime}\right|^{k} d x\right)^{2} \leqq K(k) \int|y|^{k} d x \int\left|y^{\prime \prime}\right|^{k} d x
$$

the integrals being taken over the interval in question.
[We are unable to determine the best value of $K$, even when $k=4$.]
269. If $k>1$ and $y$ and $y^{\prime \prime}$ are $L^{k}$ and $L^{k^{\prime}}$ respectively, then

$$
\int_{-\infty}^{\infty} y^{\prime 2} d x<\left(\int_{-\infty}^{\infty}|y|^{k} d x\right)^{1 / k}\left(\int_{-\infty}^{\infty}\left|y^{\prime \prime}\right|^{k^{\prime}} d x\right)^{1 / k^{\prime}}
$$

unless $y$ is nul.
270. If $y^{\prime}$ is $L^{2}$, then

$$
\begin{equation*}
\int_{0}^{\infty} \frac{y^{4}}{x^{3}} d x<\frac{3}{2}\left(\int_{0}^{\infty} y^{\prime 2} d x\right)^{2} \tag{i}
\end{equation*}
$$

unless
(ii)

$$
y=\frac{x}{a x+b},
$$

where $a$ and $b$ are positive, in which case there is equality. Moregenerally, if

$$
\begin{equation*}
l>k>1, \quad r=\frac{l}{k}-1, \tag{iii}
\end{equation*}
$$

and $y^{\prime}$ is positive and $L^{k}$, then

$$
\begin{equation*}
\int_{0}^{\infty} \frac{y^{l}}{x^{l-r}} d x<K\left(\int_{0}^{\infty} y^{\prime k} d x\right)^{l / k} \tag{iv}
\end{equation*}
$$

where
(v)

$$
K=\frac{1}{l-r-1}\left[\frac{r \Gamma(l / r)}{\Gamma(1 / r) \Gamma\{(l-1) / r\}}\right],
$$

unless

$$
\begin{equation*}
y=\frac{x}{\left(a x^{r}+b\right)^{1 / r}} . \tag{vi}
\end{equation*}
$$

[It is easy to prove an inequality of the type of (i) but with a less favourable constant. Thus, if we denote the integrals in (i) by $K$ and $J$ respectively, we have

$$
y^{2}=\left(\int_{0}^{x} y^{\prime} d t\right)^{2} \leqq J x
$$

and so

$$
K=\int_{0}^{\infty} \frac{y^{4}}{x^{3}} d x \leqq J \int_{0}^{\infty} \frac{y^{2}}{x^{2}} d x<4 J^{2}
$$

by Theorem 253. We know no elementary proof of the full result. For details of the variational proof, which is much more difficult than that of Theorem 260, we must refer to Bliss (3).]
271. If $\alpha>1$, and $\mathfrak{G}(f)$ is the geometric mean of $f$ over $(0, x)$, then

$$
\int_{0}^{\infty} x^{\alpha-1} \mathfrak{G}^{\alpha}(f) d x<\frac{1}{\alpha-1}\left\{\Gamma\left(\frac{\alpha}{\alpha-1}\right)\right\}^{1-\alpha}\left(\int_{0}^{\infty} f d x\right)^{\alpha},
$$

unless

$$
f \equiv C \exp \left(-B x^{a-1}\right)
$$

[See Hardy and Littlewood (7). The limiting case $\alpha=1$ corresponds to Theorem 335.]
272. In the problem

$$
\int_{-\infty}^{\infty} y^{2} d x \text { maximum, } \int_{-\infty}^{\infty} x^{2} y^{2} d x \text { and } \int_{-\infty}^{\infty} y^{\prime 2} d x \text { given', }
$$

the Euler equation is of the form

$$
y^{\prime \prime}+\left(a+b x^{2}\right) y=0
$$

and is soluble by parabolic cylinder functions. It has a solution $y=e^{-c x^{2}}$ if $b=-a^{2}=-(2 c)^{2}$.
[Thisgives the variational basisfor Weyl'sinequality (see Theorem 226).]

## CHAPTER VIII

## SOME THEOREMS CONCERNING BILINEAR AND MULTILINEAR FORMS

8.1. Introduction. In this chapter we prove a number of general theorems concerning the maxima of bilinear and multilinear forms. In the early part of the chapter we consider forms in $n$ sets of variables, but suppose the variables and coefficients positive. Later, we abandon this restriction, but suppose that $n=2$; and most of the latter part of the chapter is occupied by the proof of an important theorem of M. Riesz concerning bilinear forms with complex variables and coefficients.

## POSITIVE MULTILINEAR FORMS

### 8.2. An inequality for multilinear forms with positive

 variables and coefficients. We suppose that$$
x_{i}, y_{j}, \ldots, z_{k}
$$

are $n$ sets of variables, $i, j, \ldots, k$ running from $-\infty$ to $\infty$; and that

$$
\Sigma, \sum_{i}, \Sigma_{i}^{\prime}
$$

indicate respectively summation with respect to all suffixes, summation with respect to $i$ only, and summation with respect to $j, \ldots, k$ (all suffixes except $i$ ). The sum

$$
S=\Sigma a_{i j \ldots k} x_{i} y_{j} \ldots z_{k}
$$

is called a multilinear form in the variables $x, y, \ldots, z$. When $n$ is 1,2 , or 3 , the form is said to be linear, bilinear, or trilinear.

If the series are absolutely convergent, then

$$
S=\sum_{i} x_{i} \sum_{i}^{\prime} a_{i j \ldots k} y_{j} \ldots z_{k}=\sum_{j} y_{j}{\underset{j}{\prime}}_{\Sigma^{\prime}} a_{i j \ldots k} x_{i} \ldots z_{k}=\ldots
$$

273. Suppose that

$$
\begin{equation*}
0<\alpha \leqq 1, \quad 0<\beta \leqq 1, \quad \cdots, \quad 0<\gamma \leqq 1 \tag{8.2.1}
\end{equation*}
$$

and

$$
\begin{equation*}
0 \leqq \frac{\alpha+\beta+\ldots+\gamma-1}{n-1} \leqq \operatorname{Min}(\alpha, \beta, \ldots, \gamma) ; \tag{8.2.2}
\end{equation*}
$$

that $\bar{\alpha}, \bar{\beta}, \ldots, \bar{\gamma}$ are defined by
(8.2.3) $\alpha-\bar{\alpha}=\beta-\bar{\beta}=\ldots=\gamma-\bar{\gamma}=\frac{\alpha+\beta+\ldots+\gamma-1}{n-1}$
(so that $0 \leqq \bar{\alpha} \leqq \alpha, \ldots, 0 \leqq \bar{\gamma} \leqq \gamma$ ), and that $\rho, \sigma, \ldots, \tau$ are positive and
(8.2.4) $\quad \bar{\alpha} \rho+\bar{\beta} \sigma+\ldots+\bar{\gamma} \tau=1$.

Suppose further that

$$
\begin{align*}
& x_{i} \geqq 0, \quad y_{j} \geqq 0, \quad \ldots, \quad z_{k} \geqq 0, \quad a_{i j \ldots k} \geqq 0 \\
& \sum_{i} x_{i}^{1 / \alpha} \leqq X, \quad \sum_{j} y_{j}^{1 / \beta} \leqq Y, \quad \ldots, \quad \sum_{k}^{1 z_{k}}{ }^{1 / \gamma} \leqq Z \tag{8.2.5}
\end{align*}
$$

and
(8.2.6) $\quad \sum_{i}^{\prime} a_{i j \ldots k}^{\mathrm{p}}=A_{i} \leqq A, \quad{ }_{j}^{\prime} a_{i j \ldots k}^{\sigma}=B_{j} \leqq B$,

$$
\ldots, \quad \Sigma_{k}^{\prime} a_{i j \ldots k}^{\tau}=C_{k} \leqq C
$$

Then

$$
S=\Sigma a_{i j \ldots k} x_{i} y_{j} \ldots z_{k} \leqq A^{\bar{\alpha}} B^{\bar{\beta}} \ldots C^{\bar{\gamma}} X^{\alpha} Y^{\beta} \ldots Z^{\gamma}
$$

We have in fact

$$
1-\bar{\alpha}-\bar{\beta}-\ldots-\bar{\gamma}=\frac{\alpha+\beta+\ldots+\gamma-1}{n-1}=\alpha-\bar{\alpha}=\beta-\bar{\beta}=\ldots
$$

and so, by Theorem 11, a

$$
\begin{aligned}
& S=\Sigma\left(a^{\rho} x^{1 / \alpha}\right)^{\bar{\alpha}}\left(a^{\sigma} y^{1 / \beta}\right)^{\bar{\beta}} \ldots\left(a^{\tau} z^{1 / \gamma}\right)^{\bar{\gamma}}\left(x^{1 / \alpha} y^{1 / \beta} \ldots z^{1 / \gamma}\right)^{1-\bar{\alpha}-\bar{\beta}} \cdots-\bar{\gamma} \\
& \leqq\left(\Sigma a^{\rho} x^{1 / \alpha}\right)^{\bar{\alpha}}\left(\Sigma a^{\sigma} y^{1 / \beta}\right)^{\bar{\beta}} \ldots\left(\Sigma a^{\tau} z^{1 / \gamma}\right)^{\bar{\gamma}}\left(\Sigma x^{1 / \alpha} y^{1 / \beta} \ldots z^{1 / \gamma}\right)^{1-\bar{\alpha}-\bar{\beta}-\ldots-\bar{\gamma}} \\
& =\left(\sum_{i} x^{1 / \alpha} \sum_{i}^{\prime} a^{\rho}\right)^{\bar{\alpha}} \ldots\left(\sum_{k} z^{1 / \gamma} \sum_{k}^{\prime} a^{\tau}\right)^{\bar{\gamma}}\left(\sum_{i} x^{1 / \alpha} \ldots \sum_{k} z^{1 / \gamma}\right)^{1-\bar{\alpha}-\cdots-\bar{\gamma}} \\
& \leqq(A X)^{\bar{\alpha}} \ldots(C Z)^{\bar{\gamma}}(X Y \ldots Z)^{1-\bar{\alpha}-\ldots-\bar{\gamma}} \\
& =A^{\bar{\alpha}} B^{\bar{\beta}} \ldots C^{\bar{\gamma}} X^{\alpha} Y^{\beta} \ldots Z^{\gamma} .
\end{aligned}
$$

We note some special cases.
(1) If

$$
\begin{gathered}
\alpha+\beta+\ldots+\gamma=1 \\
\bar{\alpha}=\alpha, \quad \bar{\beta}=\beta, \quad \ldots, \quad \bar{\gamma}=\gamma,
\end{gathered}
$$

then
and the statement of the theorem becomes simpler.
(2) When $n=2$, the second of the conditions (8.2.2) is satisfied automatically. If we write ${ }^{b}$
then

$$
\begin{array}{ll}
\alpha=\frac{1}{p}, & \beta=\frac{1}{q} \\
\bar{\alpha}=\frac{1}{q^{\prime}}, & \bar{\beta}=\frac{1}{p^{\prime}}
\end{array}
$$

${ }^{\text {a }}$ Extended to infinite series: we shall not usually repeat this remark.
${ }^{\text {b }}$ In the preceding chapters the letters $p$ and $q$ have been reserved for the weights of mean values. In this chapter they are not required for this purpose, and we use them as indices.

Exchanging $\rho$ and $\sigma$, and $A$ and $B$, we obtain
274. If

$$
\begin{gathered}
p \geqq 1, \quad q \geqq 1, \quad \frac{1}{p}+\frac{1}{q} \geqq 1, \quad \rho>0, \quad \sigma>0, \quad \frac{\rho}{p^{\prime}}+\frac{\sigma}{q^{\prime}}=1, \\
\sum_{i} a_{i j}{ }^{\rho} \leqq A, \quad \sum_{j} a_{i j} \leqq \leqq B, \quad \sum_{i} x_{i}^{p} \leqq X, \quad \sum_{j} y_{j}^{q} \leqq Y, \\
S=\Sigma a_{i j} x_{i} y_{j} \leqq A^{1 / p^{\prime}} B^{1 / q^{\prime}} X^{1 / p} Y^{1 / q} .
\end{gathered}
$$

then
(3) An interesting and still more special case is that in which $\rho=\sigma=1, q=p^{\prime}$.
275. If

$$
p>1, \quad q>1, \quad \frac{1}{p}+\frac{1}{q}=1
$$

$$
\sum_{i} a_{i j} \leqq A, \quad \sum_{j} a_{i j} \leqq B, \quad \sum_{i} x_{i}^{p} \leqq X, \quad \sum_{j} y_{j}^{q} \leqq Y
$$

then

$$
S=\Sigma \Sigma a_{i j} x_{i} y_{j} \leqq(B X)^{1 / p}(A Y)^{1 / q}
$$

For the case $p=q=2$, see Frobenius (1) and Schur (1).
8.3. A theorem of W.H. Young. Another specialisation of Theorem 274 leads to an inequality of W. H. Young which is very important in the theory of Fourier series.

Suppose that $\sigma=\rho>1$, so that

$$
\frac{1}{p^{\prime}}+\frac{1}{q^{\prime}}=\frac{1}{\rho}<1, \quad \frac{1}{p}+\frac{1}{q}>1
$$

and take $a_{i j}=a_{i+j}$. Then

$$
\sum_{i} a_{i j}^{\rho}=\sum_{j} a_{i j}^{\rho}=\Sigma a_{n}^{\rho}=A
$$

say, for every $j$ and $i$ respectively. Hence, if we write

$$
\begin{equation*}
z_{n}=\sum_{i+j=n} x_{i} y_{j} \tag{8.3.1}
\end{equation*}
$$

we have, by Theorem 274,

$$
\begin{equation*}
\Sigma a_{n} z_{n}=\Sigma a_{i+j} x_{i} y_{j} \leqq A^{1 / \rho} X^{1 / p} Y^{1 / q} \tag{8.3.2}
\end{equation*}
$$

Since (8.3.2) is true for all $a_{n}$ for which $\Sigma a_{n}{ }^{\rho}=A$, it follows, by Theorem 15, that $\quad \Sigma z_{n}{ }^{\rho^{\prime}} \leqq X^{\rho^{\prime} / p} Y^{\rho^{\prime} / q}$.
This must be replaced, when $\rho=1$, by $z_{n} \leqq X^{1 / p} Y^{1 / q}$.
We have thus proved (apart from the specification of the cases of equality) the following theorem of Younga.
a Young $(3,4,6)$. Young does not consider the question of equality.
276. If

$$
p>1, \quad q>1, \quad \frac{1}{p}+\frac{1}{q}>1
$$

and $z_{n}$ is defined by (8.3.1), then

$$
\begin{equation*}
\Sigma z_{n}^{\frac{p q}{p+q-p q}} \leqq\left(\Sigma x_{i}^{p}\right)^{\frac{q}{p+q-p q}}\left(\Sigma y_{j}^{q}\right)^{\frac{p}{p+q-p q}} . \tag{8.3.3}
\end{equation*}
$$

Equality can occur only if all the $x$, or all the $y$, or all the $x$ but one and all the $y$ but one, are zero.

We add a more direct proof which enables us to settle the question of equality. If we write $1 / p=1-\lambda, 1 / q=1-\mu, *$ then $\lambda>0, \mu>0, \lambda+\mu<1$, and we can enunciate the theorem in the following form.
277. If $\lambda>0, \quad \mu>0, \quad \lambda+\mu<1$, $z$ is defined by (8.3.1), and $\Im_{r}(x)$ as in $\S 2.10$, then

$$
\begin{equation*}
\Im_{1 /(1-\lambda-\mu)}(z) \leqq \Im_{1 /(1-\lambda)}(x) \mathbb{S}_{1 /(1-\mu)}(y) \tag{8.3.4}
\end{equation*}
$$

with equality only in the cases specified in Theorem 276.
Let $\nu=1-\lambda-\mu$. It follows from Theorem 11 that

$$
(\Sigma u v)^{\frac{1}{\nu}}=\left(\Sigma u^{\frac{\lambda}{\lambda+\nu}} v^{\frac{\mu}{\mu+\nu}} u^{\frac{\nu}{\lambda+\nu}} v^{\frac{\nu}{\mu+\nu}}\right)^{\frac{1}{\nu}} \leqq\left(\Sigma u^{\left.\frac{1}{\lambda+\nu}\right)^{\frac{\lambda}{\nu}}}\left(\Sigma v^{\frac{1}{\mu+\nu}}\right)^{\frac{\mu}{\nu}} \Sigma u^{\frac{1}{\lambda+\nu}} v^{\frac{1}{\mu+\nu}}\right.
$$

Applying this inequality, with

$$
u=y_{n-i}=y_{j}, \quad v=x_{i}
$$

to (8.3.1), we obtain

$$
\begin{align*}
& z_{n}^{\frac{1}{\nu}} \leqq\left(\sum_{i} x_{i}^{\frac{1}{\mu+\nu}}\right)^{\frac{\mu}{\nu}}\left(\sum_{j} y_{j}^{\left.\frac{1}{\lambda+\nu}\right)^{\frac{\lambda}{\nu}}} \sum_{i+j=n} x_{i}^{\frac{1}{\mu+\nu}} y_{j}^{\frac{1}{\lambda+\nu}}\right.  \tag{8.3.5}\\
& z_{n}^{\frac{1}{\nu}} \leqq \mathbb{S}_{\frac{1}{\mu+\nu}}^{\frac{\mu}{\nu(\mu+\nu)}}(x) \mathbb{S}_{\frac{1}{\lambda+\nu}}^{\frac{\lambda}{\nu(\lambda+\nu)}}(y) \sum_{i+j=n} x_{i}^{\frac{1}{\mu+\nu}} y_{j}^{\frac{1}{\lambda+\nu}} .
\end{align*}
$$

Hence

$$
\sum_{n} z_{n}^{\frac{1}{\nu}} \leqq \mathbb{S}_{\frac{1}{\mu+\nu}}^{\frac{\mu}{\nu(\mu+\nu)}}(x) \mathbb{S}_{\frac{1}{\lambda+\nu}}^{\frac{\lambda}{\nu(\lambda+\nu)}}(y) \sum_{n} \sum_{i+j=n} x_{i}^{\frac{1}{\mu+\nu}} y_{j}^{\frac{1}{\lambda+\nu}}
$$

Since the double sum here is equal to

$$
\sum_{i} \sum_{j} x_{i}^{\frac{1}{\mu+\nu}} y_{j}^{\frac{1}{\lambda+\nu}}=\mathbb{S}_{\frac{1}{\mu+\nu}}^{\frac{1}{\mu+\nu}}(x) \mathbb{S}_{\frac{1}{\lambda+\nu}}^{\frac{1}{\lambda+\nu}}(y)
$$

we obtain

$$
\sum_{n} z_{n}^{\frac{1}{\nu}} \leqq \Im_{\frac{1}{\mu+\nu}}^{\frac{1}{\nu}}(x) \Im_{\frac{1}{\lambda+\nu}}^{\frac{1}{\nu}}(y)
$$

and this is (8.3.4).

We can dispose of the question of equality as follows.
If not all $x$, and not all $y$, are zero, there are $n$ such that $x_{i} y_{j}>0$ for some $i, j$ for which $i+j=n$. We call the lattice point corresponding to such a pair $i, j$ a point $P_{n}$. Then (8.3.5) is true, for such an $n$, if, in the first two sums on the right, we limit ourselves to values of $i$ and $j$ corresponding to points $P_{n}$. If there is equality in (8.3.5), then, for such $i, j$ the ratios

$$
x_{i}^{\frac{1}{\mu+\bar{\nu}}}: y_{j}^{\frac{1}{1+\nu}}: x_{i}^{\frac{1}{\mu+\nu}} y_{j}^{\frac{1}{\lambda+\nu}}
$$

do not depend on $i$ and $j$, and the corresponding $x_{i}$, and the corresponding $y_{j}$, are all equal. It follows that there can, for any $n$, be only a finite number of points $P_{n}$.

Suppose that all these conditions are satisfied for a certain $n$. Then equality will still be excluded in the next inequality (8.3.6), unless the $i$ and $j$ corresponding to the $P_{n}$ exhaust all $i$ and $j$ for which $x_{i}>0$ and $y_{j}>0$. It follows that the total number of positive $x_{i}$ and $y_{j}$ is finite. There is therefore a single point for which $x_{i} y_{j}>0$ and $n=i+j$ is a minimum. For this $n$ there is a unique $P_{n}$ and, if there is equality in (8.3.6) for this $n$, then $x_{i}=0$ and $y_{j}=0$ except for the corresponding pair $i, j$.
8.4. Generalisations and analogues. Theorems 276 and 277 have many interesting specialisations, generalisations, and analogues. We state a number of these without proof.
278. If $\lambda>0, \mu>0, \ldots, \nu>0, \lambda+\mu+\ldots+\nu<1$, and
then

$$
w_{n}=\underset{i+j+\ldots+k=n}{\Sigma} x_{i} y_{j} \ldots z_{k},
$$

$$
\mathbb{S}_{1-\lambda-\mu-\ldots-\nu}(w) \leqq \mathbb{S}_{\frac{1}{1-\lambda}}(x) \mathbb{S}_{\frac{1}{1-\mu}}(y) \ldots \mathbb{S}_{\frac{1}{1-\nu}}(z),
$$

unless all numbers of one set, or all but one of every set, are zero.

$$
\begin{array}{ll}
\text { 279. If } & c_{n}=\underset{i_{1}+i_{2}+\ldots+i_{k}=n}{\Sigma} a_{i_{1}} a_{i_{2}} \ldots a_{i_{k}}, \\
\text { then } & \Sigma c_{n}{ }^{2} \leqq\left(\Sigma a_{n}^{\frac{2 k}{2 k-1}}\right)^{2 k-1},
\end{array}
$$

unless all a but one are zero.

Theorems 277-279, being 'homogeneous in $\Sigma$ ' in the sense of § 1.4, have integral analogues.
280. If $\lambda>0, \mu>0, \lambda+\mu<1$,

$$
\Im_{r}(f)=\left(\int_{-\infty}^{\infty} f^{r} d x\right)^{1 / r}
$$

and

$$
h(x)=\int_{-\infty}^{\infty} f(t) g(x-t) d t,
$$

then

$$
\mathfrak{\Im}_{\frac{1}{1-\lambda-\mu}}(h)<\mathfrak{\Im}_{\frac{1}{1-\lambda}}(f) \mathfrak{\Im}_{\frac{1}{1-\mu}}(g),
$$

unless for $g$ is nul.
281. If $\lambda>0, \mu>0, \lambda+\mu<1$,

$$
\Im_{r}(f)=\int_{0}^{\infty} f^{r} d x,
$$

and

$$
h(x)=\int_{0}^{x} f(t) g(x-t) d t,
$$

then

$$
\mathfrak{J}_{\frac{1}{1-\lambda-\mu}}(h)<\mathfrak{F}_{\frac{1}{1-\lambda}}(f) \mathfrak{J}_{\frac{1}{1-\mu}}(g),
$$

unless $f$ or $g$ is nul.
282. If $k$ is an integer and

$$
\begin{aligned}
& \phi(x)=\int_{-\infty}^{\infty} \ldots \int_{-\infty}^{\infty} f\left(x_{1}\right) f\left(x_{2}\right) \ldots f\left(x_{k-1}\right) \\
& \quad \times f\left(x-x_{1}-\ldots-x_{k-1}\right) d x_{1} d x_{2} \ldots d x_{k-1},
\end{aligned}
$$

then

$$
\int_{-\infty}^{\infty} \phi^{2}(x) d x \leqq\left(\int_{-\infty}^{\infty} f^{\frac{2 k}{2 k-1}}(x) d x\right)^{2 k-1}
$$

283. If

$$
\phi(x)=\int f\left(x_{1}\right) \ldots f\left(x_{k-1}\right) f\left(x-x_{1}-\ldots-x_{k-1}\right) d x_{1} \ldots d x_{k-1},
$$

the integration being defined by
then

$$
\begin{aligned}
x_{i} & \geqq 0, \quad \Sigma x_{i} \leqq x \\
\int_{0}^{\infty} \phi^{2}(x) d x & \leqq\left(\int_{0}^{\infty} \frac{2 k}{2 k-1}(x) d x\right)^{2 k-1}
\end{aligned}
$$

284. If $f(x)$ has the period $2 \pi$, and

$$
\begin{aligned}
& \qquad \begin{array}{l}
\phi(x)= \\
\frac{1}{(2 \pi)^{k-1}} \int_{-\pi}^{\pi} \ldots \int_{-\pi}^{\pi} f\left(x_{1}\right) f\left(x_{2}\right) \ldots f\left(x_{k-1}\right) \\
\times f\left(x-x_{1}-\ldots-x_{k-1}\right) d x_{1} d x_{2} \ldots d x_{k-1}, \\
\text { then } \quad \\
\quad \frac{1}{2 \pi} \int_{-\pi}^{\pi}\{\phi(x)\}^{2} d x \leqq\left(\frac{1}{2 \pi} \int_{-\pi}^{\pi} f^{\frac{2 k}{2 k-1}}(x) d x\right)^{2 k-1}
\end{array}
\end{aligned}
$$

8.5. Applications to Fourier series. Theorems 279 and 284 have important applications in the theory of Fourier series. Here we are concerned with functions and coefficients which are not positive, but the theorems which we have proved are sufficient for the applications.

Suppose first that $f(x)$ and $g(x)$ are complex functions of $L^{2}$, and that

$$
\sum_{-\infty}^{\infty} a_{n} e^{n i x}, \quad \sum_{-\infty}^{\infty} b_{n} e^{n i x}
$$

are their complex Fourier series. Then it is well known that

$$
\begin{equation*}
\Sigma a_{n} \bar{b}_{n}=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f(x) \overline{g(x)} d x \tag{8.5.1}
\end{equation*}
$$

(the bar denoting the conjugate). In particular, if $f(x)$ is $L^{2}$,

$$
\begin{equation*}
\Sigma\left|a_{n}\right|^{2}=\frac{1}{2 \pi} \int_{-\pi}^{\pi}|f|^{2} d x \tag{8.5.2}
\end{equation*}
$$

Conversely, if $\Sigma\left|a_{n}\right|^{2}$ is convergent, there is an $f(x)$ of $L^{2}$ which has the $a_{n}$ as its Fourier constants and satisfies (8.5.2).

These theorems( 'Parseval's Theorem' and 'the Riesz-FischerTheorem') were generalised by Young and Hausdorff. We write

$$
\begin{equation*}
\mathfrak{S}_{p}(a)=\left(\Sigma\left|a_{n}\right|^{p}\right)^{1 / p}, \quad \mathfrak{I}_{p}(f)=\left(\frac{1}{2 \pi} \int_{-\pi}^{\pi}|f(x)|^{p} d x\right)^{1 / p} ; \tag{8.5.3}
\end{equation*}
$$

so that $\Im_{p}(a)$ is $\Im_{p}(|a|)$ as defined in $\S 2.10$, and (8.5.2) may be written

$$
\begin{equation*}
\mathfrak{S}_{2}(a)=\mathfrak{I}_{2}(f) \tag{8.5.4}
\end{equation*}
$$

Young and Hausdorff proved that, if

$$
\begin{equation*}
1<p \leqq 2, \tag{8.5.5}
\end{equation*}
$$

then

$$
\begin{equation*}
\mathfrak{J}_{p^{\prime}}(f) \leqq \mathfrak{S}_{p}(a) \tag{8.5.6}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathfrak{S}_{p^{\prime}}(a) \leqq \mathfrak{I}_{p}(f) \tag{8.5.7}
\end{equation*}
$$

The limitation on $p$ is essential. The theorems were proved first by Young $(3,4,6)$ for a special sequence of values of $p$ and $p^{\prime}$, viz.

$$
\begin{equation*}
p=\frac{2 k}{2 k-1}, \quad p^{\prime}=2 k \quad(k=1,2,3, \ldots), \tag{8.5.8}
\end{equation*}
$$

and then generally by Hausdorff (2).

We confine ourselves here to the case (8.5.8), considered by Young. In this case (8.5.6) and (8.5.7) are corollaries of Theorems 279 and 284 respectively. For example, the $c_{n}$ of Theorem 279a is the Fourier constant of $\psi=f^{k}$, and so

$$
\begin{aligned}
\mathfrak{J}_{2 k}^{2 k}(f) & =\frac{1}{2 \pi} \int_{-\pi}^{\pi}|f|^{2 k} d x=\frac{1}{2 \pi} \int_{-\pi}^{\pi}|\psi|^{2} d x \\
& =\Sigma\left|c_{n}\right|^{2} \leqq \mathfrak{S}_{2 k /(2 k-1)}^{2 k}(a),
\end{aligned}
$$

which is (8.5.6). Similarly $a_{n}{ }^{k}$ is the Fourier constant of the $\phi(x)$ of Theorem 284, and so (8.5.7) may be deduced from Theorem 284.

The proof of (8.5.6) and (8.5.7) for general $p$ is decidedly more difficult: see §8.17.

It is interesting to observe (as another application of Hölder's inequality) how (8.5.7) may be deduced from (8.5.6). Write

$$
b_{n}=\left|a_{n}\right|^{p^{\prime}-1} \operatorname{sgn} a_{n}=\left|a_{n}\right|^{p^{\prime}} / \bar{a}_{n}
$$

if $a_{n} \neq 0$ and $|n| \leqq N$, and $b_{n}=0$ otherwise; and let

$$
g(x)=\Sigma b_{n} e^{n i x}
$$

Then

$$
\sum_{-N}^{N}\left|a_{n}\right|^{p^{\prime}}=\sum_{-N}^{N} a_{n} \bar{b}_{n}=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f(x) \overline{g(x)} d x,
$$

since $\bar{g}$ is a trigonometrical polynomial. Hence, using Hölder's inequality and (8.5.6), we obtain

$$
\begin{aligned}
\sum_{-N}^{N}\left|a_{n}\right|^{p^{\prime}} & \leqq \mathfrak{I}_{p}(f) \mathfrak{I}_{p^{\prime}}(g) \leqq \mathfrak{I}_{p}(f) \mathfrak{S}_{p}(b) \\
& =\mathfrak{I}_{p}(f)\left(\sum_{-N}^{N}\left|a_{n}\right|^{\left(p^{\prime}-1\right) p}\right)^{1 / p}=\mathfrak{I}_{p}(f)\left(\sum_{-N}^{N}\left|a_{n}\right|^{p^{\prime}}\right)^{1 / p}
\end{aligned}
$$

Transposing the last factor, and then making $N$ tend to infinity, we obtain (8.5.7).
8.6. The convexity theorem for positive multilinear forms. In this section we prove a simple butimportant property of multilinear forms with positive variables and coefficients. The theorem which we prove is a mere corollary of Hölder's inequality, but it is useful, and will serve as an introduction to the deeper theorem of $\S 8.13$.
285. ${ }^{\text {b }}$ Suppose that $a \geqq 0, x \geqq 0, \ldots, z \geqq 0$, and that

$$
M_{\alpha, \beta, \ldots, \gamma}
$$

is the upper bound of

$$
S=\Sigma a_{i j \ldots k} x_{i} y_{j} \ldots z_{k},
$$

for all $x, y, \ldots, z$ for which

$$
\Sigma x^{1 / \alpha} \leqq 1, \quad \Sigma y^{1 / \beta} \leqq 1, \quad \ldots, \quad \Sigma z^{1 / \gamma} \leqq 1
$$

a Now, of course, formed from complex $a$.
${ }^{\mathrm{b}}$ M. Riesz (1): Riesz has $n=2$.
then $\log M_{\alpha, \beta, \ldots, \gamma}$ is a convex function of $\alpha, \beta, \ldots, \gamma$ in the region $\alpha>0, \beta>0, \ldots, \gamma>0$.

By a convex function of $n$ variables $\alpha, \beta, \ldots, \gamma$ we mean (§3.12) a function convex along anystraight line in the space of $\alpha, \beta, \ldots, \gamma$.

We have to verify that if

$$
\begin{gathered}
t_{1} \geqq 0, \quad t_{2} \geqq 0, \quad t_{1}+t_{2}=1 \\
\alpha=t_{1} \alpha_{1}+t_{2} \alpha_{2}, \quad \beta=t_{1} \beta_{1}+t_{2} \beta_{2}, \quad \cdots, \quad \gamma=t_{1} \gamma_{1}+t_{2} \gamma_{2}
\end{gathered}
$$

then

$$
\begin{equation*}
M_{\alpha, \beta, \ldots, \gamma} \leqq M^{t_{1}}{ }_{\alpha_{1}, \beta_{1}, \ldots, \gamma_{1}} M^{t_{\alpha_{2}}, \beta_{2}, \ldots, \gamma_{2}} . \tag{8.6.1}
\end{equation*}
$$

Now

$$
\begin{aligned}
S=\Sigma a x y \ldots z= & \Sigma\left(a x^{\alpha_{1} / \alpha} y^{\beta_{1} / \beta} \ldots z^{\gamma_{1} / \gamma}\right)^{t_{1}}\left(a x^{\alpha_{2} / \alpha} y^{\beta_{2} / \beta} \ldots z^{\gamma_{2} / \gamma}\right)^{t_{2}} \\
& \leqq\left(\Sigma a x^{\alpha_{1} / \alpha} y^{\beta_{1} / \beta} \ldots z^{\gamma_{1} / \gamma}\right)^{t_{1}}\left(\Sigma a x^{\alpha_{2} / \alpha} y^{\beta_{2} / \beta} \ldots z^{\gamma_{2} / \gamma}\right)^{t_{2}}
\end{aligned}
$$

Since

$$
\Sigma\left(x^{\alpha_{1} / \alpha}\right)^{1 / \alpha_{1}}=\Sigma x^{1 / \alpha} \leqq 1, \ldots,
$$

the first sum on the right does not exceed $M_{\alpha_{1}, \beta_{1}, \ldots, \gamma_{1}}$; and similarly the second does not exceed $M_{\alpha_{2}, \beta_{2}, \ldots, \gamma_{2}}$. This proves (8.6.1).

The theorem may be extended to the closed region $\alpha \geqq 0$, $\beta \geqq 0, \ldots$, if we replace the conditions $\Sigma x^{1 / \alpha} \leqq 1, \Sigma y^{1 / \beta} \leqq 1, \ldots$ by $x \leqq 1, y \leqq 1, \ldots$ when $\alpha, \beta, \ldots$ are zero.

Suppose for example that $n=2$ and

$$
\sum_{i} a_{i j} \leqq A, \quad \sum_{j} a_{i j} \leqq B
$$

Then $M_{0,1} \leqq A$ and $M_{1,0} \leqq B$, so that $M_{\alpha, 1-\alpha} \leqq B^{\alpha} A^{1-\alpha}$ for $0<\alpha<1$. If $p>1, q=p^{\prime}$, we may take $\alpha=1 / p, 1-\alpha=1 / q$; and then we obtain

$$
M_{1 / p, 1 / q} \leqq M_{1,0}^{1 / p} M_{0,1}^{1 / q} \leqq B^{1 / p} A^{1 / q}
$$

which is equivalent to the result of Theorem 275.

## GENERALITIES CONCERNING BILINEAR FORMS

8.7. General bilinear forms. So far we have been occupied with 'positive' multilinear forms, i.e. forms whose variables and coefficients are non-negative. The most important multilinear forms are bilinear, and the remainder of this chapter, and most of the next, is concerned, from one point of view or another, with bilinear forms, which will not generally be positive.

We shall denote the form whose coefficients are $a_{i j}$ by $A$, and similarly with other letters. We change the convention of §8.1 concerning the range of the suffixes; until the end of $\S 8.12, i, j, \ldots$ run from 1 to $\infty$. We write

$$
\Im_{r}(x)=\Im_{r}(|x|)=\left(\sum_{1}^{\infty}\left|x_{i}\right|^{r}\right)^{1 / r}
$$

We also write

$$
\begin{equation*}
\sum_{i} a_{i j} x_{i}=X_{j}, \quad \sum_{j} a_{i j} y_{j}=Y_{i} \tag{8.7.1}
\end{equation*}
$$

When the form is positive

$$
\begin{equation*}
A=\sum_{i} \sum_{j} a_{i j} x_{i} y_{j}=\sum_{i} x_{i} Y_{i}=\sum_{j} y_{j} X_{j}, \tag{8.7.2}
\end{equation*}
$$

the convergence of any one of these series involving the convergence of the others and the equality of the three. The equations (8.7.2) are true also for complex $a, x, y$ when the form is finite.

We shall make repeated use of the following general theorems ${ }^{\text {a }}$.
286. Suppose that

$$
p>1, q>1, \quad \frac{1}{p}+\frac{1}{p^{\prime}}=1, \quad \frac{1}{q}+\frac{1}{q^{\prime}}=1
$$

(so that $p^{\prime}>1, q^{\prime}>1$ ) and that $a, x, y$ are real and non-negative. Then the three assertions

$$
\begin{equation*}
|A(x, y)| \leqq K \Im_{p}(x) \Im_{q}(y) \tag{8.7.3}
\end{equation*}
$$

for all $x, y$; ${ }^{\text {b }}$

$$
\begin{equation*}
\mathbb{S}_{q^{\prime}}(X) \leqq K \mathbb{S}_{p}(x) \tag{8.7.4}
\end{equation*}
$$

for all $x$;

$$
\begin{equation*}
\Im_{p^{\prime}}(Y) \leqq K \Im_{q}(y) \tag{8.7.5}
\end{equation*}
$$

for all $y$; are equivalent.
287. The three assertions
(i) there is inequality in (8.7.3), unless either $\left(x_{i}\right)$ or $\left(y_{j}\right)$ is nul;
(ii) there is inequality in (8.7.4), unless $\left(x_{i}\right)$ is nul;
(iii) there is inequality in (8.7.5), unless ( $y_{j}$ ) is nul;
are also equivalent.

[^49]288. When the forms are finite, Theorems 286 and 287 are true also for forms with complex variables and coefficients.

Theorem 286 is a simple corollary of Theorems 13 and 15. It follows from (8.7.2) and Theorem 13 that

$$
\begin{equation*}
A=\sum_{j} y_{j} X_{j} \leqq \Im_{q}(y) \Im_{q^{\prime}}(X) \tag{8.7.6}
\end{equation*}
$$

so that (8.7.4) is a sufficient condition for the truth of (8.7.3); and Theorem 15 shows that it is also a necessary condition. Hence (8.7.3) and (8.7.4) are equivalent, and similarly (8.7.3) and (8.7.5) are equivalent.

There will be inequality in (8.7.3) unless $\left(y_{j}\right)$ is nul or there is equality in (8.7.4). Hence the second assertion of Theorem 287 implies the first. If the $x_{i}$, and so the $X_{j}$, are given, we can, by Theorem 13, choose a non-nul $\left(y_{j}\right)$ so that there shall be equality in (8.7.6). Hence, if there is equality in (8.7.4) for a non-nul $\left(x_{i}\right)$, then there is equality in (8.7.3) for a non-nul $\left(x_{i}\right)$ and $\left(y_{j}\right)$. Hence the two assertions are equivalent, and similarly the first and third are equivalent. This proves Theorem 287.

Finally, the whole argument applies equally to complex $a, x, y$ when the forms are finite ${ }^{\mathrm{a}}$. We have only to use Theorem 14 instead of Theorem 13.

The most important case is that in which $q=p^{\prime}, q^{\prime}=p$, when (8.7.3), (8.7.4), and (8.7.5) become
$|A| \leqq K \Im_{p}(x) \Im_{p^{\prime}}(y), \quad \Im_{p}(X) \leqq K \Im_{p}(x), \quad \Im_{p^{\prime}}(Y) \leqq K \Im_{p^{\prime}}(y)$.

## BOUNDED BILINEAR FORMS

8.8. Definition of a bounded bilinear form. Throughout the rest of this chapter we suppose, except when there is an explicit statement to the contrary, that the variables and coefficients in the forms considered are arbitrary real or complex numbers.

We describe the aggregate of all sets $(x)$ or $x_{1}, x_{2}, \ldots$, real or complex, for which

$$
\begin{equation*}
\mathfrak{S}_{p}(x)=\Im_{p}(|x|)=\left(\Sigma\left|x_{i}\right|^{p}\right)^{1 / p} \leqq 1 \tag{8.8.1}
\end{equation*}
$$

as space $[p]$. Here $p$ is any positive number; but usually $p>1$.
a Otherwise there are difficulties concerning the mode of summation of $A$.

Similarly we describe the aggregate of sets $(x, y)$ for which

$$
\begin{equation*}
\Im_{p}(x) \leqq 1, \quad \Im_{q}(y) \leqq 1 \tag{8.8.2}
\end{equation*}
$$

as space $[p, q]$. The most important case is that in which $p=q=2$. The importance of this space was first recognised by Hilbert, and we may describe it shortly as Hilbert space.

In our general definition $p$ or $q$ may be $\infty$, if we interpret $\Im_{\infty}(x)$ as $\operatorname{Max}|x|$. Thus space $[\infty, \infty]$ is the aggregate of sets $(x, y)$ for which $|x| \leqq 1,|y| \leqq 1$.

A bilinear form

$$
\begin{equation*}
A=A(x, y)=\Sigma \Sigma a_{i j} x_{i} y_{j} \tag{8.8.3}
\end{equation*}
$$

is said to be bounded in $[p, q]$ if

$$
\begin{equation*}
\left|A_{n}(x, y)\right|=\left|\sum_{i=1}^{n} \sum_{j=1}^{n} a_{i j} x_{i} y_{j}\right| \leqq M, \tag{8.8.4}
\end{equation*}
$$

where $M$ is independent of the $x$ and $y$, and of $n$, for all points of [ $p, q]$. We call $A_{n}$ a section (Abschnitt) of $A$ : a form is bounded if its sections are bounded.

It is plain that (8.8.4) will hold for all points of $[p, q]$ if it holds whenever

$$
\begin{equation*}
\mathfrak{S}_{p}(x)=1, \quad \Im_{q}(y)=1 \tag{8.8.5}
\end{equation*}
$$

In this case (8.8.4) may be written

$$
\begin{equation*}
\left|A_{n}(x, y)\right| \leqq M \Im_{p}(x) \Im_{q}(y), \tag{8.8.6}
\end{equation*}
$$

and here both sides are homogeneous of degree 1 in $x$ and in $y$, so that the conditions (8.8.5) are immaterial. We might therefore have taken (8.8.6), with unrestricted $x, y$, in our definition of a bounded form.

So far $M$ has been any number for which (8.8.4) or (8.8.6) is true; if so, we say that $A$ is bounded by $M$, or $M$ is $a$ bound of $A$. It is natural to take $M$ to be the smallest such bound, and then we say that $M$ is the bound of $A .^{\text {a }}$
a If $M_{n}$ is the maximum of $A_{n}$, under the conditions

$$
\sum_{1}^{n}\left|x_{i}\right|^{p} \leqq 1, \quad \sum_{1}^{n}\left|y_{j}\right|^{q} \leqq 1,
$$

then plainly $M_{n} \leqq M_{n+1}$ and $M_{n}$ is bounded in $n$. Hence

$$
M=\lim _{n \rightarrow \infty} M_{n}
$$

exists, and is the smallest bound of $A$.

An important special case is that in which $p=q$ and

$$
a_{i j}=a_{j i},
$$

when $A$ is said to be symmetrical. In this case a necessary and sufficient condition that $A$ should be bounded is that the quadratic form

$$
A(x, x)=\Sigma \Sigma a_{i j} x_{\imath} x_{j}
$$

should be bounded. When we say that $A(x, x)$ is bounded, we mean, naturally, that $A(x, y)$ is bounded when the $x$ and $y$ are the same, i.e. that

$$
\left|A_{n}(x, x)\right| \leqq M
$$

for all $x$ for which $\Im_{p}(x)=1$.
In the first place, it is obvious that the condition is necessary, and that the bound of $A(x, x)$ does not exceed that of $A(x, y)$. That the condition is sufficient follows from the identity

$$
\begin{equation*}
A_{n}(x, y)=\frac{1}{4} A_{n}(x+y, x+y)-\frac{1}{4} A_{n}(x-y, x-y) . \tag{8.8.7}
\end{equation*}
$$

When $p=2$ we can go a little further; the bounds of $A(x, x)$ and $A(x, y)$ are the same. In fact, if $M$ is the bound of $A(x, x)$, then (8.8.7) gives

$$
\left|A_{n}(x, y)\right| \leqq \frac{1}{4} M \Sigma\left(|x+y|^{2}+|x-y|^{2}\right)=\frac{1}{2} M \Sigma\left(|x|^{2}+|y|^{2}\right) \leqq M .
$$

It is evident that, when the coefficients $a$ are positive, $A$ is bounded if it is bounded for non-negative $x$ and $y$, and that its bound may be defined with reference only to such $x$ and $y$. If

$$
A^{*}=\Sigma \Sigma\left|a_{i j}\right| x_{i} y_{j}
$$

is bounded, $A$ is said to be absolutely bounded.
8.9. Some properties of bounded forms in $[p, q]$. The theory of bounded forms is very important, but we cannot develop it systematically here. We prove enough to enable us to give an account of some special forms in which we shall be interested in the sequel.

We take $p>1, q>1$, and, as usual,

$$
p^{\prime}=\frac{p}{p-1}, \quad q^{\prime}=\frac{q}{q-1} .
$$

289. If $A$ has the bound $M$ in $[p, q]$, then

$$
\begin{equation*}
\Sigma_{i}\left|a_{i j}\right|^{p^{\prime}} \leqq M^{p^{\prime}}, \quad \sum_{j}\left|a_{i j}\right|^{q^{\prime}} \leqq M^{q^{\prime}} \tag{8.9.1}
\end{equation*}
$$

for every $j$ and $i$ respectively.
Take all the $x$ zero except $x_{I}, x_{I}=1$, and all the $y$ zero except $y_{1}, y_{2}, \ldots, y_{J}$. Вy (8.8.6),

$$
\left|\sum_{1}^{J} a_{I_{j}} y_{j}\right| \leqq M\left(\sum_{1}^{J}\left|y_{j}\right|^{q}\right)^{1 / q}
$$

Since this is true for all $y_{j}$ and all $J$, it follows, by Theorem 15 , that

$$
\sum_{j}\left|a_{I_{j}}\right|^{q^{\prime}} \leqq M^{q^{\prime}}
$$

This is the second of (8.9.1), and the first is proved similarly.
Thus a necessary condition for boundedness in [2, 2] is that

$$
\sum_{i}\left|a_{i j}\right|^{2}<\infty, \quad \sum_{j}\left|a_{i j}\right|^{2}<\infty
$$

for all $j$ and $i$. The condition

$$
\Sigma \Sigma\left|a_{i j}\right|^{2}<\infty
$$

is sufficient, since then

$$
|A|^{2} \leqq \sum_{i, j}\left|a_{i j}\right|^{2} \sum_{i, j}\left|x_{i}{ }^{2} y_{j}{ }^{2}\right|=\sum_{i, j}\left|a_{i j}\right|^{2} \sum_{i}\left|x_{i}\right|^{2} \sum_{j}\left|y_{j}\right|^{2} ;
$$

but this condition is by no means necessary, even when the coefficients are positive. Thus, as we shall see in §8.12,

$$
\Sigma \Sigma \frac{x_{i} y_{j}}{i+j}
$$

is bounded.
290. Any row or column of a bounded form is absolutely convergent.
For

$$
\sum_{i}\left|a_{i j} x_{i} y_{j}\right| \leqq\left|y_{j}\right|\left(\sum_{i}\left|x_{i}\right|^{p}\right)^{1 / p}\left(\sum_{i}\left|a_{i j}\right|^{p^{\prime}}\right)^{1 / p^{\prime}} \leqq M\left|y_{j}\right| \Im_{p}(x)
$$

by Theorem 289.
It is plain that, when $a_{i j} \geqq 0$, a necessary condition that $A$ should be bounded is that

$$
\begin{equation*}
\Sigma \Sigma a_{i j} x_{i} y_{j} \tag{8.9.2}
\end{equation*}
$$

should be convergent for all positive $x$ and $y$ in $[p, q]$. It is natural to ask whether this is also true for bounded forms with arbitrary real or complex coefficients, i.e. whether, when $A$ is bounded, the series (8.9.2) is necessarily convergent (for the $x$ and $y$ of $[p, q]$ ) in any of the recognised senses. The answer is affirmative: if $A$ is bounded, the series (8.9.2) is convergent (indeed uniformly) in the three standard senses, as a double
series in Pringsheim's sense or as a repeated sum by rows or columns. But the idea of the convergence of the double series is irrelevant to our present purpose (and is not very important in the general theory), and we shall not prove these theorems. See Hellinger and Toeplitz (1) for the case [2,2].
8.10. The Faltung of two forms in $\left[p, p^{\prime}\right]$. We suppose now that $q=p^{\prime}$. If $A$ and $B$ are bounded in [ $\left.p, p^{\prime}\right]$, with bounds $M$ and $N$, then, by Theorem 289,

$$
\sum_{k}\left|a_{i k}\right|^{p} \leqq M^{p}, \quad \sum_{k}\left|b_{k j}\right|^{p^{\prime}} \leqq N^{p^{\prime}}
$$

and therefore, by Theorem 13, the series

$$
\begin{equation*}
f_{i j}=\sum_{k} a_{i k} b_{k j} \tag{8.10.1}
\end{equation*}
$$

is absolutely convergent. We call

$$
F=F(A, B)=\Sigma \Sigma f_{i j} x_{i} y_{j}
$$

the Faltung of $A$ and $B$. The order of $A$ and $B$ is relevant, $F(A, B)$ and $F(B, A)$ being usually different forms.
291. If $A$ and $B$ have the bounds $M$ and $N$ in $\left[p, p^{\prime}\right]$, then $F$ is bounded in $\left[p, p^{\prime}\right]$, and its bound does not exceed $M N$.

Suppose that $m \geqq n$ and that $x_{i}=0$ for $i>n$. Then, since $A$ is bounded by $M$ in [ $p, p^{\prime}$ ], we have

$$
\left|A_{m}\right|=\left|\sum_{k=1}^{m} y_{k} \sum_{i=1}^{n} a_{i k} x_{i}\right| \leqq M
$$

for all $x$ and $y$ for which $\Im_{p}(x) \leqq 1$ and $\Im_{p^{\prime}}(y) \leqq 1$. Hence, by Theorem 15,

$$
\sum_{k=1}^{m}\left|\sum_{i=1}^{n} a_{i k} x_{i}\right|^{p} \leqq M^{p}
$$

for $m \geqq n$; and therefore

$$
\begin{equation*}
\sum_{k}\left|\sum_{i=1}^{n} a_{i k} x_{i}\right|^{p} \leqq M^{p} \tag{8.10.2}
\end{equation*}
$$

Similarly

$$
\begin{equation*}
\sum_{k}\left|\sum_{j=1}^{n} b_{k j} y_{j}\right|^{p^{\prime}} \leqq N^{p^{\prime}} \tag{8.10.3}
\end{equation*}
$$

But

$$
\begin{equation*}
\sum_{i=1}^{n} \sum_{j=1}^{n} f_{i j} x_{i} y_{j}=\sum_{i=1}^{n} \sum_{j=1}^{n} x_{i} y_{j} \sum_{k} a_{i k} b_{k j}=\sum_{k}\left(\sum_{i=1}^{n} a_{i k} x_{i}\right)\left(\sum_{j=1}^{n} b_{k j} y_{j}\right) \tag{8.10.4}
\end{equation*}
$$

From (8.10.2), (8.10.3), (8.10.4), and Theorem 13, it follows that

$$
\left|\sum_{i=1}^{n} \sum_{j=1}^{n} f_{i j} x_{i} y_{j}\right| \leqq M N
$$

which proves the theorem.
It is plain that we can define the Faltung of $A$ and $B$, whether $A$ and $B$ are bounded or not, whenever the series

$$
\sum_{k}\left|a_{i k}\right|^{p}, \quad \sum_{k}\left|b_{k j}\right|^{p^{\prime}}
$$

are convergent.
8.11. Some special theorems on forms in [2, 2].a In this section we confine ourselves to the classical case $p=q=2$, and suppose the variables and coefficients real (though not generally positive). We suppose then that $A$ is a real form, and denote by

$$
\begin{equation*}
A^{\prime}=\Sigma \Sigma a_{j i} x_{i} y_{j} \tag{8.11.1}
\end{equation*}
$$

the form obtained from $A$ by exchanging suffixes in $a_{i j}$.
If

$$
\begin{equation*}
\sum_{k} a_{i k}^{2}<\infty \tag{8.11.2}
\end{equation*}
$$

for all $i$, the series

$$
\begin{equation*}
c_{i j}=\sum_{k} a_{i k} a_{j k} \tag{8.11.3}
\end{equation*}
$$

are absolutely convergent, and, by (8.10.1),

$$
\begin{equation*}
C(x, y)=\Sigma \Sigma c_{i j} x_{i} y_{j} \tag{8.11.4}
\end{equation*}
$$

is the Faltung $F\left(A, A^{\prime}\right)$ of $A$ and $A^{\prime}$. In particular $C(x, x)$ is a quadratic form whose section $C_{n}$ is given, after (8.10.4), by

$$
\begin{equation*}
C_{n}(x, x)=\sum_{k}\left(\sum_{i=1}^{n} a_{i k} x_{i}\right)^{2} \tag{8.11.5}
\end{equation*}
$$

We write

$$
C(x, x)=N(A)
$$

and call $C(x, x)$ the norm of $A$. When $A$ satisfies (8.11.2), we say that the norm of $A$ exists. The existence of $N(A)$ is, after Theorem 289, a necessary condition for the boundedness of $A$.

If $A$ is bounded, with bound $M$, then, by Theorem $289, N(A)$

[^50]is bounded, and its bound $P$ does not exceed $M^{2}$. On the other hand, whenever $N(A)$ exists,
\[

$$
\begin{aligned}
\left|A_{n}(x, y)\right| & =\left|\sum_{j=1}^{n} y_{j} \sum_{i=1}^{n} a_{i j} x_{i}\right| \leqq\left(\sum_{j=1}^{n} y_{j}^{2}\right)^{\frac{1}{2}}\left\{\sum_{j=1}^{n}\left(\sum_{i=1}^{n} a_{i j} x_{i}\right)^{2}\right\}^{\frac{1}{2}} \\
& \leqq\left(\sum_{j=1}^{n} y_{j}^{2}\right)^{\frac{1}{2}}\left\{C_{n}(x, x)\right\}^{\frac{1}{2}}
\end{aligned}
$$
\]

Hence, if $N(A)$ is bounded by $P, A$ is bounded by $P^{\frac{1}{2}}$.
Collecting our results, we obtain
292. A necessary and sufficient condition that a real form $A$ should be bounded in $[2,2]$ is that the norm $N(A)$ of $A$ should exist and be bounded. If $M$ is the bound of $A$, and $P$ that of $N(A)$, then

$$
P=M^{2}
$$

A useful corollary is
293. If $A, B, \ldots$ is a finite set of forms whose norms exist, and

$$
H(x, x)=N(A)+N(B)+\ldots
$$

is bounded, with bound $P$, then $A, B, \ldots$ are bounded, with bounds which do not exceed $P^{\frac{1}{2}}$.

In fact, if $N_{n}(A), \ldots$ are sections of $N(A), \ldots$, then $N_{n}(A), \ldots$ are non-negative ${ }^{\text {a }}$, by (8.11.5), and

$$
H_{n}(x, x)=N_{n}(A)+N_{n}(B)+\ldots
$$

Hence $N_{n}(A), \ldots$ are bounded by $P^{\frac{1}{2}}$.
8.12. Application to Hilbert's forms. We now apply Theorem 293 to two very important special forms first studied by Hilbert.
294. The forms

$$
A=\Sigma \Sigma \frac{x_{i} y_{j}}{i+j-1}, \quad B=\Sigma \Sigma^{\prime} \frac{x_{i} y_{j}}{i-j}
$$

where $i, j=1,2, \ldots$ and the dash implies the omission of the terms in which $i=j$, are bounded in real space [2, 2], with bounds not exceeding $\pi$.

[^51]It is plain that each form satisfies condition (8.11.2). We write

$$
N(A)=\Sigma \Sigma c_{i j} x_{i} x_{j}, \quad N(B)=\Sigma \Sigma d_{i j} x_{i} y_{j}
$$

and calculate $c_{i j}+d_{i j}$.
If $i=j$, we have

$$
\begin{equation*}
c_{i i}+d_{i i}=\sum_{k=1}^{\infty} \frac{1}{(i+k-1)^{2}}+\sum_{k=1}^{\infty} \frac{1}{(i-k)^{2}}=\sum_{-\infty}^{\infty} \frac{1}{(i-k)^{2}}=\frac{1}{3} \pi^{2} . \tag{8.12.1}
\end{equation*}
$$

If $i \neq j$, then

$$
\begin{aligned}
c_{i j}+d_{i j} & =\sum_{k=1}^{\infty} \frac{1}{(i+k-1)(j+k-1)}+\sum_{k=1}^{\infty} \frac{1}{(i-k)(j-k)} \\
& =\sum_{k=-\infty}^{\infty} \frac{1}{(i-k)(j-k)}=\frac{1}{i-j} \sum_{k=-\infty}^{\infty}\left(\frac{1}{j-k}-\frac{1}{i-k}\right)
\end{aligned}
$$

the dash here excluding the values $k=i$ and $k=j$. If $K$ is greater than both $|i|$ and $|j|$, then

$$
\begin{aligned}
& \underset{k=-K}{\sum^{\prime}}\left(\frac{1}{j-k}-\frac{1}{i-k}\right) \\
& \quad=\frac{2}{i-j}+\left(\frac{1}{j-K}+\cdots+\frac{1}{j+K}-\frac{1}{i-K}-\cdots-\frac{1}{i+K}\right)
\end{aligned}
$$

with two series unbroken except for the omission of terms with denominator 0 , and the bracket tends to zero when $K \rightarrow \infty$. ${ }^{\text {a }}$ Hence

$$
\begin{equation*}
c_{i j}+d_{i j}=\frac{2}{(i-j)^{2}} \quad(i \neq j) . \tag{8.12.2}
\end{equation*}
$$

From (8.12.1) and (8.12.2) it follows that

$$
\begin{equation*}
N(A)+N(B)=\frac{\pi^{2}}{3} \Sigma x_{i}{ }^{2}+2 \Sigma \Sigma^{\prime} \frac{x_{i} x_{j}}{(i-j)^{2}} . \tag{8.12.3}
\end{equation*}
$$

The first form here has the bound $\frac{1}{3} \pi^{2}$; and, since

$$
\sum_{i}^{\prime} \frac{1}{(i-j)^{2}}=\sum_{j}^{\prime} \frac{1}{(i-j)^{2}}<\frac{1}{3} \pi^{2},
$$

the second satisfies the conditions of Theorem 275, and has the bound $\frac{2}{3} \pi^{2}$. Hence $N(A)+N(B)$ has the bound $\pi^{2}$, and Theorem 294 follows from Theorem 293. ${ }^{\text {b }}$

[^52]That $A$ is bounded can be proved more simply: we give a number of proofs in Ch. IX.
$A$ is absolutely bounded (§8.8), since its coefficients are positive. It is important to observe that this is not true of $B$. To prove this, it is enough to prove that

$$
\Sigma \Sigma^{\prime} \frac{x_{i} y_{j}}{|i-j|}=\infty
$$

for a positive set $(x, y)$ for which $\Sigma x_{i}{ }^{2}$ and $\Sigma y_{j}{ }^{2}$ are convergent. We take

$$
x_{i}=i^{-\frac{1}{2}}(\log i)^{-1}(i>1), \quad y_{j}=j^{-\frac{1}{2}}(\log j)^{-1} \quad(j>1)
$$

and $x_{1}=x_{2}, y_{1}=y_{2}$. Then

$$
\begin{aligned}
\Sigma \Sigma^{\prime} \frac{x_{i} y_{j}}{|i-j|} & \geqq \sum_{j=1}^{\infty} \sum_{k=1}^{\infty} k^{-1} x_{j+k} y_{j} \geqq \sum_{k=1}^{\infty} k^{-1} \sum_{l=k+1}^{\infty} x_{l} y_{l} \\
& =\sum_{k=1}^{\infty} \frac{1}{k} \sum_{l=k+1}^{\infty} \frac{1}{l(\log l)^{2}} \geqq \sum_{k=1}^{\infty} \frac{1}{k} \int_{k+1}^{\infty} \frac{d u}{u(\log u)^{2}}=\sum_{k=1}^{\infty} \frac{1}{k \log (k+1)},
\end{aligned}
$$

and this series is divergent.
We shall see in Ch. IX that $A$ is bounded in [ $\left.p, p^{\prime}\right] . B$ is also bounded in $\left[p, p^{\prime}\right]$, but the proof is much more difficult: see M. Riesz (1,2), Titchmarsh (2, 3).

THE THEOREM OF M. RIESZ
8.13. The convexity theorem for bilinear forms with complex variables and coefficients. We prove next a very important theorem due to M. Riesza. This, like Theorem 285, asserts the convexity of $\log M_{\alpha, \beta}$, where $M_{\alpha, \beta}$ is the upper bound of a form of the type $A$; but in Riesz's theorem the form is bilinear, the $a, x, y$ are general complex numbers, and convexity is proved only in a restricted domain of $\alpha$ and $\beta$.

It is essential to Riesz's argument that $M_{\alpha, \beta}$ should be an attained maximum and not merely an upper bound; and we therefore consider a finite bilinear form

$$
\begin{equation*}
A=\sum_{i=1}^{m} \sum_{j=1}^{n} a_{i j} x_{i} y_{j} \tag{8.13.1}
\end{equation*}
$$

295. Suppose that $M_{\alpha, \beta}$ is the maximum of $A$ for

$$
\begin{equation*}
\sum_{1}^{m}\left|x_{i}\right|^{1 / \alpha} \leqq 1, \quad \sum_{1}^{n}\left|y_{j}\right|^{1 / \beta} \leqq 1 \tag{8.13.2}
\end{equation*}
$$

it being understood that, if $\alpha=0$ or $\beta=0$, these inequalities are

[^53]replaced by $\left|x_{i}\right| \leqq 1$ or $\left|y_{j}\right| \leqq 1$. Then $\log M_{\alpha, \beta}$ is convex in the triangle
\[

$$
\begin{equation*}
0 \leqq \alpha \leqq 1, \quad 0 \leqq \beta \leqq 1, \quad \alpha+\beta \leqq 1 . \tag{8.13.3}
\end{equation*}
$$

\]

We have to prove that, if $\left(\alpha_{1}, \beta_{1}\right)$ and $\left(\alpha_{2}, \beta_{2}\right)$ are two points of the triangle (8.13.3), $0<t<1$, and

$$
\begin{equation*}
\alpha=\alpha_{1} t+\alpha_{2}(1-t), \quad \beta=\beta_{1} t+\beta_{2}(1-t) \tag{8.13.4}
\end{equation*}
$$

then

$$
\begin{equation*}
M_{\alpha, \beta} \leqq M_{\alpha_{1}, \beta_{1}}^{t} M_{\alpha_{2}, \beta_{2}}^{1-t} \tag{8.13.5}
\end{equation*}
$$

After Theorem 88, it is enough to prove that (8.13.5) is true, when $\alpha$ and $\beta$ are given, for some $t$ for which $0<t<1$.a
We define $p, q, p^{\prime}, q^{\prime}$ by

$$
\begin{equation*}
\alpha=\frac{1}{p}, \quad \beta=\frac{1}{q}, \quad \frac{1}{p}+\frac{1}{p^{\prime}}=1, \quad \frac{1}{q}+\frac{1}{q^{\prime}}=1 . \tag{8.13.6}
\end{equation*}
$$

There is no finite system ( $p, q$ ) corresponding to the points $(0,1)$ and ( 1,0 ) of the triangle (8.13.3). These points are particularly important, but we may disregard them in the proof, and include them in the result, on grounds of continuity.

We may then write (8.13.2) in the form

$$
\begin{equation*}
\mathfrak{S}_{p}(x) \leqq 1, \quad \Im_{q}(\dot{y}) \leqq 1 ; \tag{8.13.7}
\end{equation*}
$$

and the inequalities (8.13.3) are equivalent to either of

$$
\begin{equation*}
q^{\prime} \geqq p \geqq 1 \tag{8.13.8}
\end{equation*}
$$

and

$$
\begin{equation*}
p^{\prime} \geqq q \geqq 1 . \tag{8.13.9}
\end{equation*}
$$

We shall also write, as in §8.7,

$$
\begin{equation*}
X_{j}=X_{j}(x)=\sum_{i} a_{i j} x_{i}, \quad Y_{i}=Y_{i}(y)=\sum_{j} a_{i j} y_{j} \tag{8.13.10}
\end{equation*}
$$

so that

$$
\begin{equation*}
A=\sum_{j} X_{j} y_{j}=\sum_{i} x_{i} Y_{i} \tag{8.13.11}
\end{equation*}
$$

or simply

$$
\begin{equation*}
A=\Sigma X y=\Sigma x Y . \tag{8.13.12}
\end{equation*}
$$

a We take for granted the continuity of $M_{a, \beta} . M_{a, \beta}$ is a maximum over the space defined by (8.13.2), a real space of $2 m+2 n$ dimensions (when we separate the real and imaginary parts of the $x$ and the $y$ ), which varies continuously with $\alpha$ and $\beta$; and the proof that this maximum is continuous, though tiresome in detail, belongs to the elements of analysis.

Theorem 286 enables us to give another definition of $M_{\alpha, \beta}$ which is more convenient for our present purpose. It is plain that $M_{\alpha, \beta}$ attains its maximum for a set ( $x, y$ ) in which
(8.13.13)

$$
\mathfrak{S}_{p}(x)=1, \quad \mathbb{S}_{q}(y)=1
$$

and $M_{\alpha, \beta}$ is the least number $K$ satisfying

$$
\begin{equation*}
|A| \leqq K \Im_{p}(x) \Im_{q}(y) \tag{8.13.14}
\end{equation*}
$$

for all such $(x, y)$. Since both sides are homogeneous of degree 1 in $x$ and $y$, the restrictions (8.13.13) are now irrelevant, and $M_{\alpha, \beta}$ may be defined as the least $K$ satisfying (8.13.14) for all $(x, y) .{ }^{\text {a }}$ By Theorem 286, this is also the least $K$ satisfying

$$
\begin{equation*}
\mathbb{S}_{q^{\prime}}(X) \leqq K \mathbb{S}_{p}(x) \tag{8.13.15}
\end{equation*}
$$

for all $x$, or
(8.13.16)

$$
\mathbb{S}_{p^{\prime}}(Y) \leqq K \mathbb{S}_{q}(y)
$$

for all $y$. We may therefore define $M_{\alpha, \beta}$ by

$$
\begin{equation*}
M_{\alpha, \beta}=\operatorname{Max} \frac{\mathbb{S}_{q^{\prime}}(X)}{\mathbb{S}_{p}(x)}=\operatorname{Max} \frac{\mathbb{S}_{p^{\prime}}(Y)}{\mathbb{S}_{q}(y)} \tag{8.13.17}
\end{equation*}
$$

the maxima being taken for all non-nul sets $x$ or $y$.
8.14. Further properties of a maximal set $(x, y)$. Suppose that $\left(x^{*}, y^{*}\right)$ is a set of $(x, y)$, subject to (8.13.7), for which $|A|$ attains its maximum, and that $X^{*}, Y^{*}$ are the corresponding values of $X, Y$. It is obvious (as we have observed already) that

$$
\begin{equation*}
\mathbb{S}_{p}\left(x^{*}\right)=1, \quad \mathbb{S}_{q}\left(y^{*}\right)=1 \tag{8.14.1}
\end{equation*}
$$

Also, as in (8.7.6),

$$
\begin{equation*}
|A| \leqq \Im_{q^{\prime}}(X) \Im_{q}(y), \quad|A| \leqq \Im_{p^{\prime}}(Y) \Im_{p}(x) . \tag{8.14.2}
\end{equation*}
$$

There must be equality in each of (8.14.2) when $x, y$ have the values $x^{*}, y^{*}$; for otherwise we could increase $|A|$ by leaving the $x, X$ unaltered and changing the $y$, or by leaving the $y, Y$ unaltered and changing the $x$. Hence

$$
M_{\alpha, \beta}=\left|A\left(x^{*}, y^{*}\right)\right|=\Im_{q^{\prime}}\left(X^{*}\right) \Im_{q}\left(y^{*}\right)=\Im_{p^{\prime}}\left(Y^{*}\right) \Im_{p}\left(x^{*}\right)
$$

Further, by Theorem 14,
or

$$
\begin{equation*}
\left|X_{j}^{*}\right|^{q^{\prime}}=\omega^{q^{\prime}}\left|y_{j}^{*}\right|^{q}, \tag{8.14.3}
\end{equation*}
$$

a This is merely a repetition of an argument used already in $\S 8.8$.
where $\omega$ is positive and independent of $j$; and

$$
\arg X_{j}^{*} y_{j}^{*}
$$

is independent of $j$. Hence

$$
M_{\alpha, \beta}=\left|A\left(x^{*}, y^{*}\right)\right|=\left|\Sigma X^{*} y^{*}\right|=\Sigma\left|X^{*} y^{*}\right|=\omega \Sigma\left|y^{*}\right|^{q}=\omega .
$$

Substituting in (8.14.3), and adding the corresponding result for $Y_{i}{ }^{*}$, we obtain

$$
\begin{equation*}
\left|X_{j}^{*}\right|=M_{\alpha, \beta}\left|y_{j}^{*}\right| q-1, \quad\left|Y_{i}^{*}\right|=M_{\alpha, \beta}\left|x_{i}^{*}\right|^{p-1} . \tag{8.14.4}
\end{equation*}
$$

8.15. Proof of Theorem 295. In what follows we suppose $(x, y)$ a maximal set (for the indices $\alpha, \beta$ ), omitting the asterisks. We write $p_{1}=1 / \alpha_{1}$, and so on, and $M, M_{1}, M_{2}$ for $M_{\alpha, \beta}, M_{\alpha_{1}, \beta_{1}}$, $M_{\alpha_{2}, \beta_{2}}$. Our use of $p, p_{1}, \ldots$ excludes the points $(0,1),(1,0)$ of the triangle, but, as we remarked in §8.13, this will not impair the proof.

By (8.14.4)

$$
\begin{aligned}
M ভ_{(p-1) p_{1}^{\prime}}^{p-1}(x) & =\left(M^{p_{1}^{\prime}} \Sigma\left|x_{i}\right|^{(p-1) p_{1}}\right)^{1 p_{1} p_{1}^{\prime}} \\
& =\left(\Sigma\left|Y_{i}\right|^{\left.\left.p_{1}\right)^{\prime}\right)^{1 / p_{1}^{\prime}}=ভ_{p_{p_{1}^{\prime}}}(Y) .} .\right.
\end{aligned}
$$

Comparing this with (8.13.17), we obtain

$$
\begin{equation*}
M \Im_{(p-1) p_{1}^{\prime}}^{p-1}(x) \leqq M_{1} \Im_{q_{1}}(y) . \tag{8.15.1}
\end{equation*}
$$

Similarly

$$
\begin{equation*}
M \mathbb{S}_{(q-1) q z^{\prime}}^{q-1}(y) \leqq M_{2} \Im_{p_{2}}(x) . \tag{8.15.2}
\end{equation*}
$$

Hence, if $0<t<1$, we have

$$
\begin{equation*}
M \Im_{(p-1) p_{1}^{\prime}}^{(p-1) t}(x) \Im_{(q-1) q_{2}^{2}}^{(q-1)}(1, t)(y) \leqq M_{1}^{t} M_{2}^{1-t} \mathbb{S}_{q_{1}}^{t}(y) \Im_{p_{2}}^{1-t}(x) . \tag{8.15.3}
\end{equation*}
$$

Let us assume provisionally that there is a $t$ between 0 and 1 which satisfies

$$
\begin{equation*}
\frac{1}{p}=\frac{t}{p_{1}}+\frac{1-t}{p_{2}}, \quad \frac{1}{q}=\frac{t}{q_{1}}+\frac{1-t}{q_{2}}, \tag{8.15.4}
\end{equation*}
$$

that is to say the equations (8.13.4), and that

$$
\begin{equation*}
\mathfrak{S}_{p_{2}}^{1-t}(x) \leqq \mathbb{S}_{(p-1) p_{1}^{\prime}}^{(p-1) t}(x), \quad \Im_{q_{1}}^{t}(y) \leqq \mathbb{S}_{(q-1) q_{2}^{2}}^{(q-1)(1,-t)}(y) . \tag{8.15.5}
\end{equation*}
$$

Then (8.15.3) and (8.15.5) will give

$$
M \leqq M_{1}^{t} M_{2}^{1-t},
$$

from which the theorem will follow.

It remains to justify the assumptions expressed by (8.15.4) and (8.15.5). Let us assume further that there are numbers $\mu$ and $v$ such that

$$
\begin{equation*}
0<\mu \leqq 1, \quad 0<\nu \leqq 1 \tag{8.15.6}
\end{equation*}
$$

$$
\begin{equation*}
p_{2}=(p-1) p_{1}^{\prime} \mu+p(1-\mu), \quad q_{1}=(q-1) q_{2}^{\prime} \nu+q(1-\nu) . \tag{8.15.7}
\end{equation*}
$$

By Theorem 18, a $r \log \mathfrak{S}_{r}(x)=\log \mathbb{S}_{r}^{r}(x)$ is a convex function of $r$; and the $x$, being a maximal set, satisfy (8.13.13). Hence

$$
\begin{equation*}
\mathfrak{S}_{p_{2}}^{p_{2}}(x) \leqq \mathbb{S}_{(p-1) p_{1}^{\prime}}^{(p-1) p_{1^{\prime}}{ }^{\mu}}(x) \mathfrak{S}_{p}^{p(1-\mu)}(x)=\mathbb{S}_{(p-1) p_{1}^{\prime}}^{(p-1) p_{1^{\prime}} \mu}(x), \tag{8.15.8}
\end{equation*}
$$

and similarly

$$
\begin{equation*}
\Im_{q_{1}}^{q_{1}}(y) \leqq \mathbb{S}_{(q-1) q_{2}^{\prime}}^{(q-1) q_{2}^{\prime} \nu}(y) \Im_{q}^{q(1-\nu)}(y)=\Im_{(q-1) q_{2}^{\prime}}^{(q-1) q_{2}^{\prime} \nu}(y) \tag{8.15.9}
\end{equation*}
$$

If finally

$$
\begin{equation*}
\frac{p_{1}^{\prime} \mu}{p_{2}}=\frac{t}{1-t}, \quad \frac{q_{2}^{\prime} v}{q_{1}}=\frac{1-t}{t} \tag{8.15.10}
\end{equation*}
$$

then (8.15.8) and (8.15.9) will be equivalent to (8.15.5).
In order to complete the proof, it is necessary to show that (8.15.4), (8.15.6), (8.15.7), and (8.15.10) are consistent. These conditions contain six equations to be satisfied by the four numbers $p, q, \mu, \nu$, and two inequalities. The first equation (8.15.10) gives

$$
\left(p_{1}^{\prime}-1\right) \mu+1=\frac{\left(p_{1}^{\prime}-1\right) p_{2}}{p_{1}^{\prime}} \frac{t}{1-t}+1=\frac{p_{2}}{p_{1}} \frac{t}{1-t}+1
$$

and

$$
p_{2}+p_{1}^{\prime} \mu=p_{2}\left(1+\frac{t}{1-t}\right)=\frac{p_{2}}{1-t}
$$

and the first equation (8.15.7) gives

$$
\frac{1}{p}=\frac{\left(p_{1}^{\prime}-1\right) \mu+1}{p_{2}+p_{1}^{\prime} \mu}=\frac{p_{2} t+p_{1}(1-t)}{p_{1} p_{2}}=\frac{t}{p_{1}}+\frac{1-t}{p_{2}}
$$

which agrees with (8.15.4).
A similar argument applies to the equations involving $q$, so that (8.15.4) is a consequence of (8.15.7) and (8.15.10). Given $p_{1}, q_{1}, p_{2}, q_{2}$, and $t$, we can find $\mu, \nu$ from (8.15.10) and $p, q$ from (8.15.7), and these numbers will satisfy the six equations.

[^54]It remains only to examine the inequalities (8.15.6). If $\mu$ and $\nu$ satisfy (8.15.10), and $0<t<1$, the inequalities are equivalent to

$$
\begin{equation*}
\frac{q_{1}}{q_{2}^{\prime}} \leqq \frac{t}{1-t} \leqq \frac{p_{1}^{\prime}}{p_{2}} . \tag{8.15.11}
\end{equation*}
$$

Since ( $\alpha_{1}, \beta_{1}$ ) and ( $\alpha_{2}, \beta_{2}$ ) lie in the triangle (8.13.3), we have, by (8.13.8) and (8.13.9), and $a$ fortiori

$$
q_{1} \leqq p_{1}^{\prime}, \quad q_{2}^{\prime} \geqq p_{2}
$$

$$
\frac{q_{1}}{q_{2}^{\prime}} \leqq \frac{p_{1}^{\prime}}{p_{2}} .
$$

We can therefore choose $t$ so as to satisfy (8.15.11), and then all our conditions are satisfied.

It will be observed that it is only in the last paragraph that we use the essential inequality $\alpha+\beta \geqq 1$. When the form is positive, this inequality is irrelevant; $\log M_{\alpha, \beta}$ is then, by Theorem 285, convex in the whole of the positive quadrant of $(\alpha, \beta)$.
8.16. Applications of the theorem of M. Riesz. (i) Theorem 295 is easily transformed into another theorem of very different appearance.
296. Suppose that

$$
\begin{equation*}
X_{j}(x)=\sum_{i=1}^{m} a_{i j} x_{i} \quad(j=1,2, \ldots, n) \tag{8.16.1}
\end{equation*}
$$

and that $M_{\alpha, \gamma}^{*}$ is the maximum of
for

$$
\begin{aligned}
& \left(\sum_{1}^{n}\left|X_{j}\right|^{1 / \gamma}\right)^{\gamma} \\
& \sum_{1}^{m}\left|x_{i}\right|^{1 / \alpha} \leqq 1
\end{aligned}
$$

Then $\log M_{\alpha, \gamma}^{*}$ is convex in the triangle

$$
\begin{equation*}
0 \leqq \gamma \leqq \alpha \leqq 1 \tag{8.16.2}
\end{equation*}
$$

In fact, by (8.13.17),

$$
\begin{gathered}
M_{\alpha, \beta}=\operatorname{Max} \frac{\mathfrak{S}_{q^{\prime}}(X)}{\mathbb{S}_{p}(x)} \\
M_{\alpha, \gamma}^{*}=\operatorname{Max} \frac{\mathbb{S}_{1 / \gamma}(X)}{\mathbb{S}_{1 / \alpha}(x)}=\operatorname{Max} \frac{\mathbb{S}_{1 / \gamma}(X)}{\mathbb{S}_{p}(x)} \\
M_{\alpha, \beta}=M_{\alpha, \gamma}^{*}
\end{gathered}
$$

while
Hence
if $\gamma=1 / q^{\prime}=1-\beta$; and then the conditions (8.16.2) are equivalent to (8.13.8) or to (8.13.3). Thus $M_{\alpha, \gamma}^{*}$ is a convex function of $(\alpha, 1-\gamma)$, or, what is the same thing, of $(\alpha, \gamma)$.
(ii) 297. Suppose that $X_{j}$ is defined by (8.16.1); that

$$
\begin{equation*}
\sum_{1}^{n}\left|X_{j}\right|^{2} \leqq \sum_{1}^{m}\left|x_{i}\right|^{2} \tag{8.16.3}
\end{equation*}
$$

for all $x$; and that

$$
\begin{equation*}
\mathbf{l} \leqq p \leqq 2 . \tag{8.16.4}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mathfrak{S}_{p^{\prime}}(X) \leqq \mathrm{m}^{(2-p) / p} \mathfrak{S}_{p}(x), \tag{8.16.5}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathfrak{m}=\operatorname{Max}\left|a_{i j}\right| . \tag{8.16.6}
\end{equation*}
$$

To deduce Theorem 297 from Theorem 296, we write $\alpha=1 / p$, as before, and consider the line from $\left(\frac{1}{2}, \frac{1}{2}\right)$ to $(1,0)$ in the plane $(\alpha, \gamma)$. This line lies entirely in the triangle (8.16.2); and hence, by Theorem 296,

$$
M_{\alpha, 1-\alpha}^{*} \leqq\left(M_{\hat{i}, \frac{1}{2}}^{*}\right)^{2(1-\alpha)}\left(M_{1,0}^{*}\right)^{2 \alpha-1}
$$

for $\frac{1}{2}<\alpha<1$. It is plain from (8.16.3) that $M_{2,1}^{*} \leqq 1$; and

$$
M_{1,0}^{*}=\frac{\operatorname{Max}|X|}{\Sigma|x|} \leqq m
$$

Hence

$$
M_{\alpha, 1-\alpha}^{*} \leqq \mathfrak{m}^{2 \alpha-1}=\mathfrak{m}^{(2-p) / p},
$$

which is equivalent to (8.16.5).
The condition (8.16.3) is certainly satisfied (with equality) if (8.16.1) is a 'unitary' substitution, i.e. a substitution which leaves $\Sigma|x|^{2}$ unaltered ${ }^{\text {a }}$. This case of the theorem was found by F. Riesz (4), and the general theorem by M. Riesz (1).
8.17. Applications to Fourier series. Out of many other important applications of Riesz's theorems, we select an application of Theorem 297 to the proof of Hausdorff's theorems ${ }^{\text {b }}$.

[^55](i) Suppose that $m$ is an odd integer and
\[

$$
\begin{gathered}
f_{m}(\theta)=\sum_{-\frac{1}{2} m}^{\ddagger m} e^{\mu \theta i} x_{\mu},{ }^{\mathrm{a}} \\
X_{\nu}=\sum_{-\frac{1}{2} m} a_{\mu \nu} x_{\mu}=\frac{1}{\sqrt{ } m} f_{m}\left(\frac{2 \pi \nu}{m}\right), \\
a_{\mu \nu}=m^{-\frac{1}{2}} e^{2 \mu \nu \pi i!m} .
\end{gathered}
$$
\]

The substitution is unitary, so that

$$
\Sigma\left|X_{\nu}\right|^{2}=\Sigma\left|x_{\mu}\right|^{2} .
$$

Also $\mathrm{m}=m^{-\frac{1}{2}}$. Hence, by Theorem 297,

$$
\begin{equation*}
\left(\frac{1}{m} \Sigma\left|f_{m}\left(\frac{2 \pi \nu}{m}\right)\right|^{p^{\prime}}\right)^{1 / p^{\prime}} \leqq\left(\Sigma\left|x_{\mu}\right|^{p}\right)^{1 / p} \tag{8.17.1}
\end{equation*}
$$

The left-hand side being an approximation to $\Im_{p^{\prime}}(f)$, we may deduce Hausdorff's theorem (8.5.6) by passages to the limit ${ }^{\text {b }}$.
(ii) If $m$ is again an odd integer,

$$
x_{\mu}=f_{m}\left(\frac{2 \pi \mu}{m}\right)=\sum_{r=-\frac{1}{2} m}^{\frac{1}{2} m} a_{r} e^{2 r \mu \pi i / m}
$$

and

$$
X_{\nu}=\frac{1}{\sqrt{ } m} \sum_{-\frac{k}{2} m}^{\sum_{m} e^{-2 \nu \mu \pi i / m} x_{\mu}, ~}
$$

then simple calculations show that
and

$$
\begin{gathered}
X_{\nu}=m^{\frac{1}{2}} a_{\nu}, \\
\Sigma\left|X_{\nu}\right|^{2}=\Sigma\left|x_{\mu}\right|^{2}
\end{gathered}
$$

as before. In this case Theorem 297 gives

$$
\left(\Sigma\left|a_{\nu}\right|^{p^{\prime}}\right)^{1 / p^{\prime}} \leqq\left(\frac{1}{m} \Sigma\left|f_{m}\left(\frac{2 \pi \mu}{m}\right)\right|^{p}\right)^{1 / p} ;
$$

and Hausdorff's theorem (8.5.7) follows by appropriate passages to the limit.

We can also, as we showed in §8.5, deduce the second theorem from the first.
a We now write $\mu, \nu$ for $i, j$, and extend the summations over the range

$$
\begin{aligned}
& \text { b If } f(\theta) \text { is a polynomial } \\
& \sum_{-\frac{1}{2} M}^{\frac{1}{2} M} x_{\mu} e^{\mu \theta i},
\end{aligned}
$$

then $f_{m}(\theta)=f(\theta)$ for $m \geqq M$, and the theorem for $f(\theta)$ follows immediately from (8.17.1). The extension to an arbitrary $f(\theta)$ depends on the theory of 'strong convergence'.

## MISCELLANEOUS THEOREMS AND EXAMPLES

298. If $p>1$, and $a(x, y)$ is measurable and positive, then the three assertions
(i) $\int_{0}^{\infty} \int_{0}^{\infty} a(x, y) f(x) g(y) d x d y \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{p^{\prime}} d y\right)^{1 / p^{\prime}}$,
for all non-negative $f, g$;
(ii)

$$
\int_{0}^{\infty} d y\left(\int_{0}^{\infty} a(x, y) f(x) d x\right)^{p} \leqq K^{p} \int_{0}^{\infty} f^{p} d x,
$$

for all non-negative $f$;

$$
\begin{equation*}
\int_{0}^{\infty} d x\left(\int_{0}^{\infty} a(x, y) g(y) d y\right)^{p^{\prime}} \leqq K^{p^{\prime}} \int_{0}^{\infty} g^{p^{\prime}} d y \tag{iii}
\end{equation*}
$$

for all non-negative $g$, are equivalent. The assertions 'there is inequality in (i) unless $f$ or $g$ is nul', 'there is inequality in (ii) unless $f$ is nul', 'there is inequality in (iii) unless $g$ is nul', are also equivalent.
[Analogue of Theorems 286 and 287, with $q=p^{\prime}$. Thereis a more general form with both $p$ and $q$ arbitrary.]
299. The forms

$$
A=\Sigma \Sigma \frac{x_{i} y_{j}}{i+j-1+\lambda}, \quad B=\Sigma \Sigma^{\prime} \frac{x_{i} y_{j}}{i-j+\lambda}
$$

where $\lambda>0$, and the dash is required only if $\lambda$ is integral, are bounded in [2,2], and have bounds $\pi$ if $\lambda$ is integral, $\pi|\operatorname{cosec} \lambda \pi|$ if $\lambda$ is non-integral.
[Schur (1), Pólya and Szegö (1, 1, 117, 290).]
300. If $p>1$ and $A=\Sigma \Sigma a_{i j} x_{i} y_{j}$ has a bound $M$ in [ $\left.p, p^{\prime}\right]$, and

$$
h_{i j}=\int f_{\imath}(t) g_{j}(t) d t
$$

where

$$
\begin{gathered}
\int\left|f_{i}\right|^{p} d t \leqq \mu^{p}, \quad \int\left|g_{j}\right|^{p^{\prime}} d t \leqq \nu^{p^{\prime}}, \\
A^{*}=\Sigma \Sigma a_{i j} h_{i j} x_{i} y_{j}
\end{gathered}
$$

then
has a bound $M_{\mu \nu}$.

$$
\text { [For } \quad \begin{aligned}
\left|A^{*}\right| & =\left|\int\left\{\Sigma \Sigma a_{i j} x_{i} f_{i}(t) y_{j} g_{j}(t)\right\} d t\right| \\
& \left.\leqq M \int \Sigma\left|x_{i} f_{i}(t)\right|^{p}\right\}^{1 / p}\left\{\Sigma\left|y_{j} g_{j}(t)\right|^{p^{\prime}}\right\}^{1 / p^{\prime}} d t \\
& \leqq M\left\{\int \Sigma\left|x_{i}\right|^{p}\left|f_{i}(t)\right|^{p} d t t^{1 / p}\left\{\int \Sigma\left|y_{j}\right|^{p^{\prime}}\left|g_{j}(t)\right|^{p^{\prime}} d t\right\}^{1 / p^{\prime}}\right. \\
& \leqq M \mu \nu\left(\Sigma\left|x_{i}\right|^{p}\right)^{1 / p}\left(\Sigma\left|y_{j}\right|^{p^{p}}\right)^{1 / p^{\prime}} .
\end{aligned}
$$

For the case $p=p^{\prime}=2$, see $\operatorname{Schur}$ (1).]
301. $\Sigma \Sigma^{\prime} \frac{(i j)^{\frac{1}{2}(\mu-1)}}{i^{\mu}-j^{\mu}} x_{i} y_{j}$ is bounded in [2, 2].
302. $\Sigma \Sigma^{\prime} \frac{\sin (i-j) \theta}{i-j} x_{i} y_{j}$ is bounded in [2, 2], for any real $\theta$. If $0 \leqq \theta \leqq \pi$,
the bound does not exceed $\operatorname{Max}(\theta, \pi-\theta)$.
[For the last two theorems see Schur (1).]
303. If $a_{n}$ is the Fourier sine coefficient of an odd bounded function, or the Fourier cosine coefficient of an even bounded function, then the forms
are bounded in [2, 2].

$$
\Sigma \Sigma a_{i+j} x_{i} y_{j}, \quad \Sigma \Sigma a_{i-j} x_{i} y_{j}
$$

[Toeplitz (1). Suppose that $i$ and $j$ run from 1 to $n$; that $x$ and $y$ are real and

$$
\Sigma x_{i}{ }^{2}=\Sigma y_{j}{ }^{2}=1 ;
$$

that $\quad X=\Sigma x_{i} \cos i \theta, \quad X^{\prime}=\Sigma x_{i} \sin i \theta, \quad Y=\Sigma y_{j} \cos j \theta, \quad Y^{\prime}=\Sigma y_{j} \sin j \theta$, and that (for example)

$$
a_{n}=\frac{2}{\pi} \int_{0}^{\pi} f(\theta) \sin n \theta d \theta,
$$

where $|f| \leqq M$. A simple calculation gives

$$
\sum_{1}^{n} \sum_{1}^{n} a_{i-j} x_{i} y_{j}=\frac{2}{\pi} \int_{0}^{\pi}\left(X^{\prime} Y-X Y^{\prime}\right) f(\theta) d \theta
$$

Since

$$
\left|\int_{0}^{\pi} X Y^{\prime} f(\theta) d \theta\right| \leqq \frac{1}{2} M \int_{0}^{\pi}\left(X^{2}+Y^{\prime 2}\right) d \theta=\frac{1}{4} M \pi\left(\Sigma x^{2}+\Sigma y^{2}\right)=\frac{1}{2} M \pi
$$

we find the upper bound $2 M$ for $A$. Similarly in the other cases.
If for example $f(\theta)$ is odd and equal to $\frac{1}{2}(\pi-\theta)$ for $0<\theta<\pi$, then $M=\frac{1}{2} \pi$ and $a_{n}=n^{-1}$. We thus obtain the result of Theorem 294 concerning B.]
304. If (i) $\Sigma \Sigma a_{i j} x_{i} y_{j}$ is bounded in [p,q], (ii) $k>1, l>1$, and (iii) $\left(u_{i}\right)$, $\left(v_{j}\right)$ are given sets for which $\mathbb{S}_{p k^{\prime}}(u)<\infty, \mathbb{S}_{q l^{\prime}}(v)<\infty$, then
is bounded in $[p k, q l]$.

$$
A=\Sigma \Sigma a_{i j} u_{i} v_{j} x_{i} y_{j}
$$

$$
\text { [For } \quad \begin{aligned}
|A| & \leqq M\left(\Sigma|u x|^{p}\right)^{1 / p}\left(\Sigma|v y|^{q}\right)^{1 / q} \\
& \left.\leqq M\left(\Sigma\left|u^{p k^{\prime}}\right|\right)^{1 / p k^{\prime}}\left(\Sigma\left|v^{q l^{\prime}}\right|\right)^{1 / l^{\prime}}\left(\Sigma|x|^{p k}\right)^{1 / p k}\left(\Sigma|y|^{q l}\right)^{1 / q l} .\right]
\end{aligned}
$$

305. The form

$$
\Sigma \Sigma^{\prime} \frac{u_{i} v_{j}}{i-j} x_{i} y_{j}
$$

where $u_{i}$ and $v_{j}$ are given sets of numbers satisfying

$$
\Sigma\left|u_{i}\right|^{2} \leqq 1, \quad \Sigma\left|v_{j}\right|^{2} \leqq 1,
$$

is bounded, but not necessarily absolutely bounded, in $[\infty, \infty]$.
[Take $p=q=2, k=l=\infty$ in Theorem 304.
If the form were always absolutely bounded, then Hilbert's form $B$ of § 8.12 would be absolutely bounded, which is untrue.]
306. If

$$
\begin{gathered}
\mathfrak{M}_{2}(a) \geqq A_{2} H, \quad \mathfrak{M}_{4}(a) \leqq A_{4} H, \\
\frac{A_{2}{ }^{3}}{A_{4}{ }^{2}} H \leqq \mathfrak{M}_{1}(a) \leqq A_{4} H .
\end{gathered}
$$

[By Theorems 16 and 17. The result is required in the proof of the next theorem.]
307. If perpendiculars are drawn from the corners ( $\pm 1, \pm 1, \ldots, \pm 1$ ) of the 'unit cube' in space of $m$ dimensions, upon any linear [ $m-1$ ] through the centre of the cube, then the mean of the perpendiculars lies between two constants $A$ and $B$, independent of $m$ and the position of the [ $m-1]$.
308. If

$$
b_{j}=\left(\sum_{i}\left|a_{i j}\right|^{2}\right)^{\frac{1}{2}}, \quad c_{i}=\left(\sum_{j}\left|a_{i j}\right|^{2}\right)^{\frac{1}{2}}
$$

then

$$
P=\left(\Sigma \Sigma\left|a_{i j}\right|^{\left.\frac{1}{3}\right)^{\frac{3}{2}}} \leqq K\left(\Sigma b_{j}+\Sigma c_{i}\right)=K(B+C),\right.
$$

where $K$ is an absolute constant.
309. A necessary condition that a form $A=\Sigma \Sigma a_{i j} x_{i} y_{j}$, with real or complex coefficients, should be bounded in $[\infty, \infty]$, with bound $M$, is that, in the notation of the last theorem, $B, C, P$ should be less than $K M$.
[For the last five theorems see Littlewood (2).]
310. If

$$
\begin{gathered}
p \geqq 2, q \geqq 2, \frac{1}{p}+\frac{1}{q} \leqq \frac{1}{2}, \\
\lambda=\frac{p q}{p q-p-q}, \quad \mu=\frac{4 p q}{3 p q-2 p-2 q},
\end{gathered}
$$

and $A$ is bounded in $[p, q]$, with bound $M$, then

$$
\left(\Sigma b_{j}^{\lambda}\right)^{1 / \lambda} \leqq K M, \quad\left(\Sigma c_{i}{ }^{\lambda}\right)^{1 / \lambda} \leqq K M, \quad\left(\Sigma \Sigma\left|a_{i j}\right|^{\mu}\right)^{1 / \mu} \leqq K M,
$$

where $b_{j}$ and $c_{i}$ are defined as in Theorem 308, and $K$ depends on $p$ and $q$ only.
311. If

$$
\frac{1}{2} \leqq \frac{1}{p}+\frac{1}{q}<1,
$$

but the conditions of Theorem 310 are satisfied in other respects, then

$$
\left(\Sigma b_{j}^{\lambda}\right)^{1 / \lambda} \leqq K M, \quad\left(\Sigma c_{i}\right)^{1 / \lambda} \leqq K M, \quad\left(\Sigma \Sigma\left|a_{i j}\right|^{\lambda}\right)^{1 / \lambda} \leqq K M .
$$

312. If

$$
p<2<q, \frac{1}{p}+\frac{1}{q}<1
$$

and the conditions of Theorem 310 are satisfied otherwise, then

$$
\begin{gathered}
\left(\Sigma \Sigma\left|a_{i j}\right|^{\lambda}\right)^{1 / \lambda} \leqq K M . \\
p>1, q>1, \frac{1}{p}+\frac{1}{q}<1, \\
a_{i j} \geqq 0,
\end{gathered}
$$

313. If
$A$ is bounded in $[p, q]$, with bound $M$, and

$$
\beta_{j}=\left(\sum_{i} a_{i j^{p^{\prime}}}\right)^{1 / p^{\prime}}, \quad \gamma_{i}=\left(\sum_{j} a_{i j^{q^{\prime}}}\right)^{1 / q^{\prime}},
$$

then

$$
\left(\Sigma \beta_{j}^{\lambda}\right)^{1 / \lambda} \leqq M, \quad\left(\Sigma \gamma_{i}\right)^{1 / \lambda} \leqq M, \quad\left(\Sigma \Sigma a_{i j}{ }^{\lambda}\right)^{1 / \lambda} \leqq M .
$$

[For the last four theorems see Hardy and Littlewood (13).]
314. Hilbert's forms in [ $\left.p, p^{\prime}\right]$. It will be proved in Ch. IX that the form $A$ of Theorem 294 is bounded in [ $p, p^{\prime}$ ]. The corresponding theorem for $B$ lies a good deal deeper. We have to show that

$$
\begin{equation*}
\left|\frac{\sum_{11}^{n} \sum^{\prime}}{x_{i} y_{j}} \overline{i-j}\right| \leqq K \Im_{p}(x) \Im_{p^{\prime}}(y), \tag{i}
\end{equation*}
$$

where $K=K(p)$ depends only on $p$; or, what is, after Theorem 286, the same thing, that
(ii)

$$
\sum_{i}\left|\Sigma_{j}^{\prime} \frac{y_{j}}{i-j}\right|^{p^{\prime}} \leqq K^{p^{\prime}} \Sigma\left|y_{j}\right|^{p^{\prime}} .
$$

It is enough, after Theorem 295, to prove (i) or (ii) for even integral values of $p^{\prime}$ (or for some subsequence of these values). This demands some special device, the most natural, from our present point of view, being that used by Titchmarsh (2).

## CHAPTER IX

## HILBERT'S INEQUALITY AND ITS ANALOGUES AND EXTENSIONS

9.1. Hilbert's double series theorem. The researches of which we give an account in this chapter originate in a remarkable bilinear form which was first studied by Hilbert, and which we have already encountered in §8.12, viz. the form

$$
\Sigma \Sigma \frac{a_{m} b_{n}}{m+n},
$$

where $m$ and $n$ run from 1 to $\infty$. Ourfirst theorem isTheorem 315 below, which we state with its integral analogue and with a complement of a type which will occur frequently in this chapter.
315. If
and

$$
p>1, \quad p^{\prime}=p /(p-1)
$$

the summations running from 1 to $\infty$, then

$$
\begin{equation*}
\Sigma \Sigma \frac{a_{m} b_{n}}{m+n}<\frac{\pi}{\sin (\pi / p)} A^{1 / p} B^{1 / p^{\prime}} \tag{9.1.1}
\end{equation*}
$$

unless $(a)$ or (b) is nul.
316. If

$$
p>1, \quad p^{\prime}=p /(p-1)
$$

and

$$
\int_{0}^{\infty} f^{p}(x) d x \leqq F, \quad \int_{0}^{\infty} g^{p^{\prime}}(y) d y \leqq G,
$$

then

$$
\begin{equation*}
\int_{0}^{\infty} \int_{0}^{\infty} \frac{f(x) g(y)}{x+y} d x d y<\frac{\pi}{\sin (\pi / p)} F^{1 / p} G^{1 / p^{\prime}}, \tag{9.1.2}
\end{equation*}
$$

unless $f \equiv 0$ or $g \equiv 0$.
317. The constant $\pi \operatorname{cosec}(\pi / p)$ is the best possible constant in each of Theorems 315 and 316.

The case $p=p^{\prime}=2$ of Theorem 315 is 'Hilbert's double series theorem', and was proved first (apart from the exact determination of the constant) by Hilbert in his lectures on integral equations. Hilbert's proof was published by Weyl (2). The determination of the constant, and the
integral analogue, are due to Schur (1), and the extensions to general $p$ to Hardy and M. Riesz: see Hardy (3). Other proofs, of the whole or of parts of the theorems, and generalisations in different directions, have been given by Fejér and F. Riesz (1), Francis and Littlewood (1), Hardy (2), Hardy, Littlewood, and Pólya (1), Mulholland (1,3), Owen (1), Pólya and Szegö (1, 1, 117, 290), Schur (1), and F. Wiener (1). A number of these generalisations will be proved or quoted later.

The inequality (9.1.1) is of the same type as the general inequality discussed in §8.2; but Theorem 315 is not included in Theorem 275, since

$$
\sum_{m} \frac{1}{m+n}
$$

is divergent. It is to be observed that $\pi \operatorname{cosec} \alpha \pi$, where $\alpha=1 / p$, is (in accordance with Theorem 295) convex for $0<\alpha<1$.
9.2. A general class of bilinear forms. We shall deduce Theorem 315 from the following more general theorem ${ }^{\text {a }}$.
318. Suppose that $p>1, p^{\prime}=p /(p-1)$, and that $K(x, y)$ has the following properties:
(i) $K$ is non-negative, and homogeneous of degree -1 :

$$
\begin{equation*}
\int_{0}^{\infty} K(x, 1) x^{-1 / p} d x=\int_{0}^{\infty} K(1, y) y^{-1 / p^{\prime}} d y=k \tag{ii}
\end{equation*}
$$

and either (iii) $K(x, 1) x^{-1 / p}$ is a strictly decreasing function of $x$, and $K(1, y) y^{-1 / p^{\prime}}$ of $y$ : or, more generally, (iii') $K(x, 1) x^{-1 / p}$ decreases from $x=1$ onwards, while the interval $(0,1)$ can be divided into two parts, $(0, \xi)$ and $(\xi, 1)$, of which one may be nul, in the first of which it decreases and in the second of which it increases; and $K(1, y) y^{-1 / p^{\prime}}$ has similar properties. Finally suppose that, when only the less stringent condition (iii') is satisfied,

$$
\begin{equation*}
K(x, x)=0 \tag{iv}
\end{equation*}
$$

Then
(a)

$$
\Sigma \Sigma K(m, n) a_{m} b_{n}<k\left(\Sigma a_{m}{ }^{p}\right)^{1 / p}\left(\Sigma b_{n}{ }^{p^{\prime}}\right)^{1 / p^{\prime}}
$$

unless (a) or (b) is nul;

$$
\begin{equation*}
\sum_{n}\left(\sum_{m} K(m, n) a_{m}\right)^{p}<k^{p} \Sigma a_{m}^{p} \tag{b}
\end{equation*}
$$

unless $(a)$ is nul;
a Hardy, Littlewood, and Pólya (1). The case $p=2$ of the theorem is due in substance to $\operatorname{Schur}(\mathbf{1})$ : Schur supposes $K(x, y)$ a decreasing function of both variables.
(c)

$$
\sum_{m}\left(\sum_{n} K(m, n) b_{n}\right)^{p^{\prime}}<k^{p^{\prime}} \Sigma b_{n}^{p^{\prime}}
$$

unless (b) is nul.
In each case the summations are from 1 to $\infty$. Theorems 286 and 287 show that the three conclusions (a), (b), (c) are equivalent.

We may elucidate the hypotheses by the following remarks.
(1) The convergence and equality of the two integrals in condition (ii) is a consequence of the convergence of either, because of the homogeneity of $K$.
(2) The words 'decreasing',.. are to be interpreted in the strict sense throughout the theorem.
(3) $\xi$ may be 0 or 1 , one of the intervals $(0, \xi)$ and $(\xi, 1)$ then disappearing.
(4) In the most important application, in which

$$
K(x, y)=\frac{1}{x+y},
$$

condition (iii) is satisfied. An interesting case in which condition (iii') is required is

$$
K(x, y)=\frac{1}{(x+y)^{1-a}|x-y|^{a}} \quad(0<\alpha<1) \cdot .^{a}
$$

In this case $K(x, 1)$ has an infinity at $x=1$. In such cases condition (iv) is required in order to exclude equal pairs ( $m, m$ ) from the summation.

It is easy to see that, if $m$ and $n$ are positive integers, and the summations are over $r=1,2, \ldots$, then
(9.2.1) $\Sigma K\left(\frac{r}{n}, 1\right)\left(\frac{r}{n}\right)^{-1 / p} \frac{1}{n}<\int_{0}^{\infty} K(x, 1) x^{-1 / p} d x=k$,
(9.2.2) $\Sigma K\left(1, \frac{r}{m}\right)\left(\frac{r}{m}\right)^{-1 / p^{\prime}} \frac{1}{m}<\int_{0}^{\infty} K(1, y) y^{-1 / p^{\prime}} d y=k$.

For, if (iii) is satisfied, then
(9.2.3) $K\left(\frac{r}{n}, 1\right)\left(\frac{r}{n}\right)^{-1 / p} \frac{1}{n}<\int_{(r-1) / n}^{r / n} K(x, 1) x^{-1 / p} d x$,
and (9.2.1) follows by summation. If only (iii') is satisfied, we use (9.2.3) for $r>n$ and for $r \leqq \xi n$, and

$$
K\left(\frac{r}{n}, 1\right)\left(\frac{r}{n}\right)^{-1 / p} \frac{1}{n}<\int_{r / n}^{(r+1) / n} K(x, 1) x^{-1 / p} d x
$$

a See Hardy, Littlewood, and Pólya (1).
for $\xi n<r<n$; and the result again follows by summation when we observe that $K(1,1)=0$. The proof of (9.2.2) is similar.

Hence

$$
\begin{aligned}
\Sigma \Sigma K(m, n) a_{m} b_{n} & =\Sigma \Sigma a_{m} K^{1 / p}\left(\frac{m}{n}\right)^{1 / p p^{\prime}} b_{n} K^{1 / p^{\prime}}\left(\frac{n}{m}\right)^{1 / p p^{\prime}} \\
& \leqq P^{1 / p} Q^{1 / p^{\prime}}
\end{aligned}
$$

where

$$
\begin{aligned}
P=\sum_{m} a_{m}^{p} & \sum_{n} K(m, n)\left(\frac{m}{n}\right)^{1 / p^{\prime}} \\
& =\sum_{m} a_{m}{ }^{p} \sum_{n} K\left(1, \frac{n}{m}\right)\left(\frac{n}{m}\right)^{-1 / p^{\prime}} \frac{1}{m}<k \Sigma a_{m}^{p},
\end{aligned}
$$

by (9.2.2), unless ( $a$ ) is nul; and similarly

$$
Q<k \Sigma b_{n}{ }^{p^{\prime}}
$$

unless $(b)$ is nul. This proves the theorem.
If we take $\quad K(x, y)=\frac{1}{x+y}$
we obtain Theorem 315. It may be shown that the $k$ of Theorem 318 is a best possible constant, but in this direction we shall not go beyond proving Theorem 317. a
9.3. The corresponding theorem for integrals. The theorem for integrals corresponding to Theorem 318 is
319. Suppose that $p>1$, that $K(x, y)$ is non-negative and homogeneous of degree -1 , and that

$$
\int_{0}^{\infty} K(x, 1) x^{-1 / p} d x=\int_{0}^{\infty} K(1, y) y^{-1 / p^{\prime}} d y=k
$$

Then (a)

$$
\int_{0}^{\infty} \int_{0}^{\infty} K(x, y) f(x) g(y) d x d y \leqq k\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{p^{\prime}} d y\right)^{1 / p^{\prime}}
$$

$$
\begin{align*}
& \int_{0}^{\infty} d y\left(\int_{0}^{\infty} K(x, y) f(x) d x\right)^{p} \leqq k^{p} \int_{0}^{\infty} f^{p} d x  \tag{b}\\
& \int_{0}^{\infty} d x\left(\int_{0}^{\infty} K(x, y) g(y) d y\right)^{p^{\prime}} \leqq k^{p^{\prime}} \int_{0}^{\infty} g^{p^{\prime}} d y
\end{align*}
$$

If $K(x, y)$ is positive, then there is inequality in (b) unless $f \equiv 0$, in (c) unless $g \equiv 0$, and in (a) unless either $f \equiv 0$ or $g \equiv 0$.
${ }^{\text {a }}$ See $\S 9.5$. The constant $k$ is (again in accordance with Theorem 295) convex in $\alpha=1 / p$.

The theorem may be proved by the method of $\S 9.2$, which naturally goes rather more simply in this case. We have

$$
\begin{aligned}
& \iint K(x, y) f(x) g(y) d x d y \\
& =\iint f(x) K^{1 / p}\left(\frac{x}{y}\right)^{1 / p p^{\prime}} g(y) K^{1 / p^{\prime}}\left(\frac{y}{x}\right)^{1 / p p^{\prime}} d x d y \leqq P^{1 / p} Q^{1 / p^{\prime}}, \\
& \text { where } \quad P=\int f^{p}(x) d x \int K(x, y)\left(\frac{x}{y}\right)^{1 / p^{\prime}} d y=k \int f^{p} d x \\
& \text { and } \quad Q=k \int g^{p^{\prime}} d y .
\end{aligned}
$$

If $K>0$, and there is equality, then

$$
\begin{equation*}
A f^{p}(x)\left(\frac{x}{y}\right)^{1 / p^{\prime}} \equiv B g^{p^{\prime}}(y)\left(\frac{y}{x}\right)^{1 / p} \tag{9.3.1}
\end{equation*}
$$

for almost all $y$. ${ }^{\text {a }}$ If we give $y$ a value for which $g(y)$ is positive and finite, and for which the equivalence holds, we see that $f^{p}(x)$ is equivalent to a function $C x^{-1}$, and this is inconsistent with the convergence of $\int f^{p} d x$. Hence either $f$ or $g$ is nul. It may also be shown that the constant is the best possible.

There is another interesting method of proof due to Schur ${ }^{\text {b }}$. We have

$$
\begin{aligned}
& \int_{0}^{\infty} f(x) d x \int_{0}^{\infty} K(x, y) g(y) d y=\int_{0}^{\infty} f(x) d x \int_{0}^{\infty} x K(x, x w) g(x w) d w \\
& =\int_{0}^{\infty} f(x) d x \int_{0}^{\infty} K(1, w) g(x w) d w=\int_{0}^{\infty} K(1, w) d w \int_{0}^{\infty} f(x) g(x w) d x
\end{aligned}
$$

(if any of the integrals is convergent). Applying Theorem 189 to the inner integral, and observing that

$$
\int g^{p^{\prime}}(x w) d x=\frac{1}{w} \int g^{p^{\prime}}(y) d y
$$

we obtain (a); and (b) and (c) are corollaries, by Theorem 191.
The case $K(x, y)=1 /(x+y)$ gives Theorem 316. We shall return later ( $\S 9.9$ ) to other applications.

[^56]9.4. Extensions of Theorems 318 and 319. (1) The following theorem is in some ways more and in others less general than Theorem 318.
320. Suppose that $K(x, y)$ is a strictly decreasing function of $x$ and $y$, and also satisfies the conditions (i) and (ii) of Theorem 318; that $\lambda_{m}>0, \mu_{n}>0$,
$$
\Lambda_{m}=\lambda_{1}+\lambda_{2}+\ldots+\lambda_{m}, \quad \mathbf{M}_{n}=\mu_{1}+\mu_{2}+\ldots+\mu_{n} ;
$$
and that $p>1$. Then
$$
\Sigma \Sigma K\left(\Lambda_{m}, \mathrm{M}_{n}\right) \lambda_{m}^{1 / p^{\prime}} \mu_{n}^{1 / p} a_{m} b_{n}<k\left(\Sigma a_{m}^{p}\right)^{1 / p}\left(\Sigma b_{n}^{p^{\prime}}\right)^{1 / p^{\prime}}
$$
unless (a) or (b) is nula.
The special case $\Lambda_{m}=m, \mathrm{M}_{n}=n$ is also a special case of Theorem 318.

We deduce Theorem 320 from Theorem 319 by a process which has many applications ${ }^{b}$. We interpret $\Lambda_{0}$ and $\mathrm{M}_{0}$ as 0 , and take, in Theorem 319,

$$
\begin{array}{ll}
f(x)=\lambda_{m}^{-1 / p} a_{m} & \left(\Lambda_{m-1} \leqq x<\Lambda_{m}\right), \\
g(y)=\mu_{n}-1 / p^{\prime} b_{n} & \left(\mathrm{M}_{n-1} \leqq y<\mathrm{M}_{n}\right) .
\end{array}
$$

If we observe that

$$
\int_{\Lambda_{m-1}}^{\Lambda_{m}} \int_{\mathrm{M}_{n-1}}^{\mathrm{M}_{n}} K(x, y) f(x) g(y) d x d y>\lambda_{m}^{1 / p^{\prime}} \mu_{n}^{1 / p} a_{m} b_{n} K\left(\Lambda_{m}, \mathrm{M}_{n}\right)
$$

unless $a_{m}=0$ or $b_{n}=0$, we obtain Theorem 320.
If $K(x, y)=1 /(x+y)$ we obtain ${ }^{\text {c }}$

- 321. $\Sigma \Sigma \frac{\lambda_{m}^{1 / p^{\prime}} \mu_{n}^{1 / p}}{\Lambda_{m}+\mathrm{M}_{n}} a_{m} b_{n}<\frac{\pi}{\sin (\pi / p)}\left(\Sigma a_{m}^{p}\right)^{1 / p}\left(\Sigma b_{n}{ }^{p^{\prime}}\right)^{1 / p^{\prime}}$
unless (a) or (b) is nul.
(2) Theorems 318 and 319 may be extended to multiple series or integrals of any order.
322.d Suppose that the $n$ numbers $p, q, \ldots, r$ satisfy

$$
p>1, \quad q>1, \quad \ldots, \quad r>1, \quad \frac{1}{p}+\frac{1}{q}+\ldots+\frac{1}{r}=1 ;
$$

[^57]that $K(x, y, \ldots, z)$ is a positive function of the $n$ variables $x, y, \ldots, z$, homogeneous of degree $-n+1$; and that
\[

$$
\begin{equation*}
\int_{0}^{\infty} \ldots \int_{0}^{\infty} K(1, y, \ldots, z) y^{-\frac{1}{q}} \ldots z^{-\frac{1}{r}} d y \ldots d z=k \tag{9.4.1}
\end{equation*}
$$

\]

Then

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{0}^{\infty} \cdots \int_{0}^{\infty} K(x, y, \ldots, z) f(x) g(y) \ldots h(z) d x d y \ldots d z \\
& \leqq k\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{q} d y\right)^{1 / q} \cdots\left(\int_{0}^{\infty} h^{r} d z\right)^{1 / r}
\end{aligned}
$$

If further

$$
y^{-1 / q} \ldots z^{-1 / r} K(1, y, \ldots, z), \quad x^{-1 / p} \ldots z^{-1 / r} K(x, 1, \ldots, z), \ldots
$$

are decreasing functions of all the variables which they involve, then
$\Sigma \Sigma \ldots \Sigma K(m, n, \ldots, s) a_{m} b_{n} \ldots c_{s} \leqq k\left(\Sigma a_{m}{ }^{p}\right)^{1 / p}\left(\Sigma b_{n}^{q}\right)^{1 / q} \ldots\left(\Sigma c_{s}^{r}\right)^{1 / r}$.
In virtue of the homogeneity of $K$, the convergence of (9.4.1) implies the convergence and equality of all the $n$ integrals of the same type.

Theorem 322 may be proved bystraightforward generalisations of the proofs of Theorems 318 and 319.
9.5. Best possible constants : proof of Theorem 317. We have still to prove Theorem 317, which asserts that the constant $\pi \operatorname{cosec}(\pi / p)$ of Theorems 315 and 316 is 'the best possible', that is to say that the inequalities asserted by the theorems would be false, for some $a_{m}, b_{n}$ or $f(x), g(y)$, if $\pi \operatorname{cosec}(\pi / p)$ were replaced by any smaller number. The method which we use illustrates an important general principle and may be used in the proof of many theorems of this 'negative' character.

We take $\quad a_{m}=m^{-(1+\epsilon) / p}, \quad b_{n}=n^{-(1+\epsilon) / p^{\prime}}$,
where $\epsilon$ is small and positive; we may suppose that $\epsilon<p^{\prime} / 2 p$. We denote by $O(1)$ a number which may depend upon $p$ and $\epsilon$, but is bounded when $p$ is fixed and $\epsilon \rightarrow 0$; and by $o(1)$ a number which satisfies these conditions and tends to zero with $\epsilon$. Then

$$
\frac{1}{\epsilon}=\int_{1}^{\infty} x^{-1-\epsilon} d x<\sum_{1}^{\infty} m^{-1-\epsilon}<1+\int_{1}^{\infty} x^{-1-\epsilon} d x=1+\frac{1}{\epsilon},
$$

and so
(9.5.1) $\quad \Sigma a_{m}^{p}=\Sigma m^{-1-\epsilon}=\frac{1}{\epsilon}+O(1), \quad \Sigma b_{n}^{p^{\prime}}=\frac{1}{\epsilon}+O(1)$.

Also

$$
\begin{aligned}
\Sigma \Sigma \frac{a_{m} b_{n}}{m+n}>\int_{1}^{\infty} \int_{1}^{\infty} x^{-(1+\epsilon) / p} & y^{-(1+\epsilon) / p^{\prime}} \frac{d x d y}{x+y} \\
& =\int_{1}^{\infty} x^{-1-\epsilon} d x \int_{1 / x}^{\infty} u^{-(1+\epsilon) / p^{\prime}} \frac{d u}{1+u}
\end{aligned}
$$

The error in replacing the lower limit in the inner integral by 0 is less than $x^{-\alpha} / \alpha$, where $\alpha$ is positive and independent of $\epsilon ;{ }^{\text {a }}$ and

$$
\frac{1}{\alpha} \int_{1}^{\infty} x^{-1-\alpha-\epsilon} d x<\frac{1}{\alpha^{2}}
$$

Hence
(9.5.2) $\Sigma \Sigma \frac{a_{m} b_{n}}{m+n}>\int_{1}^{\infty} x^{-1-\epsilon} d x \int_{0}^{\infty} u^{-(1+\epsilon) / p^{\prime}} \frac{d u}{1+u}+O(1)$

$$
=\frac{1}{\epsilon}\left\{\frac{\pi}{\sin (\pi / p)}+o(1)\right\}+O(1)=\frac{1}{\epsilon}\left\{\frac{\pi}{\sin (\pi / p)}+o(1)\right\}
$$

It is plain from (9.5.1) and (9.5.2) that, if $k$ is any number less than $\pi \operatorname{cosec}(\pi / p)$, then

$$
\Sigma \Sigma \frac{a_{m} b_{n}}{m+n}>k\left(\Sigma a_{m}^{p}\right)^{1 / p}\left(\Sigma b_{n}^{p^{\prime}}\right)^{1 / p^{\prime}}
$$

when $\epsilon$ is sufficiently small.
This proves that the constant in (9.1.1) is the best possible. Since (9.1.1) can be deduced from (9.1.2) as in §9.4, it follows that the constant in (9.1.2) is also the best possible. We could of course also prove this directly.

An alternative method is to take

$$
a_{m}=m^{-1 / p}, \quad b_{n}=n^{-1 / p^{\prime}}
$$

when $m \leqq \mu, n \leqq \mu$, and $a_{m}=0, b_{n}=0$ otherwise, and to make $\mu$ tend to infinity. The principle is the same in each case; we make $a_{m}$ and $b_{n}$ depend upon a parameter ( $\epsilon$ or $\mu$ ) in such a manner that the series involved tend to infinity when the parameter tends
a It is less than
where

$$
\begin{aligned}
& \int_{0}^{1 / x} u-(1+e) \left\lvert\, p^{\prime} d u=\frac{x^{-\beta}}{\beta}\right. \\
& \beta=1-\frac{1+\epsilon}{p^{\prime}}=\frac{1}{p}-\frac{\epsilon}{p^{\prime}}
\end{aligned}
$$

and we may take $\alpha=1 / 2 p$, if $\epsilon<p^{\prime} / 2 p$.
to a limit, and compare their values for values of the parameter near this limit. The method is effective in the proof of very many theorems of the type of Theorem 317. The inequalities (9.1.1) and (9.1.2) assert unattained upper bounds; except when both sides vanish, equality cannot occur; and it is for this reason that the introduction of a parameter $(\epsilon$ or $\mu)$ is necessary in the proof of the complementary theorem.
9.6. Further remarks on Hilbert's theorems ${ }^{\text {a }}$. Theorems 315 and 316 have been proved in many different ways and have very varied applications. In this and the next section we collect a number of remarks which concern both proofs and applications, and are intended to illustrate the connections between the theorems and various parts of the theory of functions.
(1) Theorem 315 may be deduced from Theorem 316 by the process which led us to Theorem 321. We define $f(x)$ and $g(y)$ by

$$
f(x)=a_{m}(m-1 \leqq x<m), \quad g(y)=b_{n}(n-1 \leqq x<n)
$$

and observe that then

$$
\int_{m-1}^{m} \int_{n-1}^{n} \frac{f(x) g(y)}{x+y} d x d y \geqq \frac{a_{m} b_{n}}{m+n}
$$

Here, however, we can go a little further, since

$$
\frac{1}{m+n-1-\alpha}+\frac{1}{m+n-1+\alpha}>\frac{2}{m+n-1}
$$

for $0<\alpha<1$, and so ${ }^{b}$

$$
\int_{m-1}^{m} \int_{n-1}^{n} \frac{d x d y}{x+y}>\frac{1}{m+n-1}
$$

If now we replace $m$ and $n$ by $m+1$ and $n+1$, we obtain a slightly sharper form of Theorem 315, viz.
323. If the conditions of Theorem 315 are satisfied, then

$$
\sum_{0}^{\infty} \sum_{0}^{\infty} \frac{a_{m} b_{n}}{m+n+1}<\frac{\pi}{\sin (\pi / p)}\left(\sum_{0}^{\infty} a_{m}^{p}\right)^{1 / p}\left(\sum_{0}^{\infty} b_{n}^{p^{\prime}}\right)^{1 / p^{\prime}} .
$$

[^58]Several other proofs of Hilbert's theorem, for example the proof of Mulholland (1), the proof of Schur given in §8.12, the proof of Fejér and F. Riesz given below, and the proof of Pólya and Szegö (1, 1, 290), also give the result in this form. The last three are limited to the case $p=2$.
(2) The proof of Fejér and F. Riesz is based upon the theory of analytic functions, and proceeds as follows. Suppose that $f(z)=\Sigma a_{n} z^{n}$ is a polynomial of degree $N$ with non-negative coefficients, not identically zero. Then, by Cauchy's Theorem,

$$
\int_{-1}^{1} f^{2}(x) d x=-i \int_{0}^{\pi} f^{2}\left(e^{i \theta}\right) e^{i \theta} d \theta
$$

and so

$$
\begin{gather*}
\int_{0}^{1} f^{2}(x) d x<\int_{-1}^{1} f^{2}(x) d x \leqq \frac{1}{2} \int_{-\pi}^{\pi}\left|f\left(e^{i \theta}\right)\right|^{2} d \theta  \tag{9.6.1}\\
\sum_{0}^{N} \sum_{0}^{N} \frac{a_{m} a_{n}}{m+n+1}<\pi \sum_{0}^{N} a_{n}^{2}
\end{gather*}
$$

If we make $N \rightarrow \infty$ we obtain Hilbert's theorem, with $a_{m}=b_{m}$ and ' $\leqq$ ' for ' $<$ '. The first restriction is unimportant, since, after $\S 8.8$, a symmetric bilinear form in [2, 2] has a bound equal to that of the corresponding quadratic form. To replace ' $\leqq$ ' by ' $<$ ' requires a refinement of the argument which we shall not discuss here.

The second inequality in (9.6.1) may be written

$$
\int_{-1}^{1}|f(x)|^{2} d x \leqq \frac{1}{2} \int_{-\pi}^{\pi}\left|f\left(e^{i \theta}\right)\right|^{2} d \theta
$$

and in this form it is valid whether the coefficients $a_{n}$ are real or complex, and has important function-theoretic applications ${ }^{\text {a }}$.
(3) Hilbert's original proof depended upon the identity

$$
\begin{equation*}
\int_{-\pi}^{\pi} t\left\{\sum_{1}^{n}(-1)^{r}\left(a_{r} \cos r t-b_{r} \sin r t\right)\right\}^{2} d t=2 \pi(S-T) \tag{9.6.2}
\end{equation*}
$$

where

$$
S=\sum_{1}^{n} \sum_{1}^{n} \frac{a_{r} b_{s}}{r+s}, \quad T=\sum_{1}^{n} \sum_{1}^{n} \frac{a_{r} b_{s}}{r-s}
$$

(the dash implying that pairs $r, s$ for which $r=s$ are omitted). From this it follows that

$$
\begin{equation*}
2 \pi|S-T| \leqq \pi \int_{-\pi}^{\pi}\left\{\sum_{1}^{n}(-1)^{r}\left(a_{r} \cos r t-b_{r} \sin r t\right)\right\}^{2} d t=\pi^{2} \sum_{1}^{n}\left(a_{r}^{2}+b_{r}^{2}\right) \tag{9.6.3}
\end{equation*}
$$

a See Fejér and F. Riesz (1). The inequality is actually true (and in the stricter form with ' $<$ ') for any $f(z)$, except $f(z)=0$, for which $\Sigma\left|a_{n}\right|^{2}$ is convergent; and this is a corollary of Hilbert's theorem, if this theorem is proved in some other way.

If $a_{r}=b_{r}, T$ disappears, and we obtain

$$
\begin{equation*}
\sum_{1}^{n} \sum_{1}^{n} \frac{a_{r} a_{s}}{r+s} \leqq \pi \sum_{1}^{n} a_{r}^{2} ; \tag{9.6.4}
\end{equation*}
$$

and from (9.6.4), and the remark of § 8.8 quoted under (2), we deduce

$$
\begin{equation*}
\sum_{1}^{n} \sum_{1}^{n} \frac{a_{r} b_{s}}{r+s} \leqq \pi\left(\sum_{1}^{n} a_{r}^{2}\right)^{\frac{1}{2}}\left(\sum_{1}^{n} b_{r}^{2}\right)^{\frac{1}{2}} \leqq \frac{1}{2} \pi\left(\sum_{1}^{n} a_{r}^{2}+\sum_{1}^{n} b_{r}^{2}\right) . \tag{9.6.5}
\end{equation*}
$$

From (9.6.3) and (9.6.5) it follows that

$$
|T|=\left|\sum_{1}^{n} \sum_{1}^{n} \frac{a_{r} b_{s}}{r-s}\right| \leqq \pi\left(\sum_{1}^{n} a_{r}^{2}+\sum_{1}^{n} b_{r}^{2}\right) ;
$$

and hence, on grounds of homogeneity, that

$$
\left|\sum_{1}^{n} \sum_{1}^{n} \frac{a_{r} b_{s}}{r-s}\right| \leqq 2 \pi\left(\sum_{1}^{n} a_{r}^{2}\right)^{\frac{1}{2}}\left(\sum_{1}^{n} b_{r}^{2}\right)^{\frac{1}{2}} .
$$

This gives the second result of Theorem 294, except that the constant $2 \pi$ is not the best possible constant.
9.7. Applications of Hilbert's theorems. (1) As an application of Hilbert's theorem to the theory of analytic functions we select the following. Suppose that $f(z)$ is regular in $|z|<1$ and belongs to the 'complex Lebesgue class $L$ ', i.e. that

$$
\frac{1}{2 \pi} \int_{-\pi}^{\pi}\left|f\left(r e^{i \theta}\right)\right| d \theta
$$

is bounded for $r<1$. If $f(z)$ is 'wurzelfrei', i.e. has no zeros in $|z|<1$, then

$$
f(z)=\Sigma c_{n} z^{n}=g^{2}(z)=\left(\Sigma a_{n} z^{n}\right)^{2}
$$

where $g(z)$ also is regular in $|z|<1$. Since $\int\left|g\left(r e^{i \theta}\right)\right|^{2} d \theta$ is bounded, $\Sigma\left|a_{n}\right|^{2}$ is convergent, and therefore, by Theorem 323,

$$
\Sigma \frac{\left|a_{m}\right|\left|a_{n}\right|}{m+n+1}
$$

is convergent. A fortiori

$$
\sum_{\nu} \frac{\left|c_{\nu}\right|}{\nu+1}=\left.\sum_{\nu} \frac{1}{v+1}\right|_{m+n=\nu} ^{\sum} a_{m} a_{n}
$$

is convergent.
It is fairly easy, by a method which is familiar in this part of the theory of analytic functions, to extend the conclusion to general $f$ (not necessarily 'wurzelfrei') ${ }^{\text {a }}$. We thus obtain the

[^59]theorem: if $f(z)$ belongs to $L$ in $|z|<1$, then its integrated power series is absolutely convergent for $|z|=1 .{ }^{\text {a }}$
(2) As an application of Hilbert's series theorem to the theory of functions of a real variable, we prove
324.b If $f(x)$ is real, $L^{2}$, and not nul, in ( 0,1 ), and
\[

$$
\begin{gathered}
a_{n}=\int_{0}^{1} x^{n} f(x) d x \quad(n=0,1,2, \ldots), \\
\Sigma a_{n}^{2}<\pi \int_{0}^{1} f^{2}(x) d x .
\end{gathered}
$$
\]

The constant is the best possible.
We may plainly suppose $f(x) \geqq 0$. Then, if $\left(b_{n}\right)$ is any nonnegative and non-nul sequence,

$$
\begin{aligned}
\Sigma a_{n} b_{n} & =\Sigma b_{n} \int_{0}^{1} x^{n} f(x) d x=\int_{0}^{1}\left(\Sigma b_{n} x^{n}\right) f(x) d x \\
\left(\Sigma a_{n} b_{n}\right)^{2} & \leqq \int_{0}^{1}\left(\Sigma b_{n} x^{n}\right)^{2} d x \int_{0}^{1} f^{2}(x) d x \\
& =\Sigma \Sigma \frac{b_{m} b_{n}}{m+n+1} \cdot \int_{0}^{1} f^{2}(x) d x<\pi \Sigma b_{n}^{2} \cdot \int_{0}^{1} f^{2}(x) d x
\end{aligned}
$$

by Hilbert's theorem. The result now follows from Theorem 15.
To prove the constant $\pi$ the best possible, consider

$$
f(x)=(1-x)^{\epsilon-\frac{1}{2}},
$$

and make $\epsilon$ tend to zero.
The integrals $a_{n}$ are called the moments of $f(x)$ in $(0,1)$ and are important in many theories.
Here we have deduced Theorem 324 from Theorem 323 (with $p=2$ ) and Theorem 15 (the converse of Hölder's inequality). We can, if we please, reverse the argument, deducing Theorem 323 (with $p=2$ ) from Theorem 324 and Theorem 191 (the integral analogue of Theorem 15). Suppose that $g(x)=\Sigma b_{n} x^{n}$, where $b_{n}$ is

[^60]non-negative and not always 0 , and that $f(x)$ is any non-negative and non-nul function. Then
\[

$$
\begin{aligned}
& \int_{0}^{1} f g d x=\int_{0}^{1}\left(\Sigma b_{n} x^{n}\right) f d x=\Sigma b_{n} \int_{0}^{1} x^{n} f d x=\Sigma a_{n} b_{n} \\
& \left(\int_{0}^{1} f g d x\right)^{2}=\left(\Sigma a_{n} b_{n}\right)^{2} \leqq \Sigma a_{n}^{2} \Sigma b_{n}^{2}<\pi \Sigma b_{n}^{2} \int_{0}^{1} f^{2} d x
\end{aligned}
$$
\]

by Theorem 324. Since this is true for all $f$, it follows, by Theorem 191, that

$$
\int_{0}^{1} g^{2} d x<\pi \Sigma b_{n}^{2}
$$

and this is equivalent to Theorem 323.
It is plain that when two inequalities, each involving a constant factor, are 'reciprocal' in this sense, each being deducible from the other in this way by the converse of Hölder's inequality, then one constant must be best possible if the other is. We shall meet with another application of this principle later ( $\S 9.10(1))$.
(3) As a corollary of Theorem 316 (with $p=2$ ), we prove
325.a Suppose that $a_{n} \geqq 0$ and that the summations run from 0 to $\infty$, and that

$$
\begin{equation*}
A(x)=\Sigma a_{n} x^{n}, \quad A^{*}(x)=\Sigma \frac{a_{n} x^{n}}{n!} \tag{9.7.1}
\end{equation*}
$$

Then

$$
\begin{gather*}
\Sigma \Sigma \frac{a_{m} a_{n}}{m+n+1} \leqq \pi \Sigma \Sigma \frac{(m+n)!}{m!n!} \frac{a_{m} a_{n}}{2^{m+n+1}},  \tag{9.7.2}\\
\int_{0}^{1} A^{2}(x) d x \leqq \pi \int_{0}^{\infty}\left\{e^{-x} A^{*}(x)\right\}^{2} d x . \tag{9.7.3}
\end{gather*}
$$

It may be verified at once by expansion and term-by-term integration that (9.7.2) and (9.7.3) are equivalent.

To prove (9.7.3) we observe that

$$
A(x)=\int_{0}^{\infty} e^{-t} A^{*}(x t) d t=\frac{1}{x} \int_{0}^{\infty} e^{-u / x} A^{*}(u) d u
$$

and so

$$
\begin{aligned}
\int_{0}^{1} A^{2}(x) d x=\int_{0}^{1} \frac{d x}{x^{2}}\left(\int_{0}^{\infty} e^{-u / x} A^{*}(u) d u\right)^{2} & =\int_{1}^{\infty} d y\left(\int_{0}^{\infty} e^{-u y} A^{*}(u) d u\right)^{2} \\
& =\int_{0}^{\infty} d w\left(\int_{0}^{\infty} e^{-u w} \alpha(u) d u\right)^{2}
\end{aligned}
$$

where

$$
\alpha(u)=e^{-u} A^{*}(u) .
$$

a Widder (1), Hardy (9).

This is

$$
\begin{aligned}
& \int_{0}^{\infty} d w \int_{0}^{\infty} e^{-u w} \alpha(u) d u \int_{0}^{\infty} e^{-v w} \alpha(v) d v=\int_{0}^{\infty} \int_{0}^{\infty} \frac{\alpha(u) \alpha(v)}{u+v} d u d v \\
& \leqq \pi \int_{0}^{\infty} \alpha^{2}(u) d u=\pi \int_{0}^{\infty}\left\{e^{-u} A^{*}(u)\right\}^{2} d u
\end{aligned}
$$

by Theorem 316.
It is easy to see that the constant $\pi$ is the best possible. The relations between the functions (9.7.1) are important in the theory of divergent series, particularly in connection with singularities of analytic functions.
9.8. Hardy's inequality. The two theorems which we discuss next were discovered in the course of attempts to simplify the proofs then known of Hilbert's theorems ${ }^{\text {a }}$.

We might require only an imperfect form of Theorem 315: the double series is convergent whenever $\Sigma a_{n}^{p}$ and $\Sigma b_{n}^{p^{\prime}}$ are convergent. It would then be natural to argue as follows. We divide the double series into two parts $S_{1}, S_{2}$ by the diagonal $m=n$, and consider the part $S_{1}$ in which $m \leqq n$. Then

$$
S_{1}=\sum_{m \leqq n} \frac{a_{m} b_{n}}{m+n} \leqq \sum_{m \leqq n} \frac{a_{m} b_{n}}{n}=\Sigma \frac{A_{n}}{n} b_{n}
$$

where

$$
A_{n}=a_{1}+a_{2}+\ldots+a_{n}
$$

Since $\Sigma b_{n}^{p^{\prime}}$ is convergent, the last series is convergent whenever $\Sigma n^{-p} A_{n}^{p}$ is convergent, and, to prove the convergence of $S_{1}$, it is enough to prove that the convergence of the last series is a consequence of that of $\Sigma a_{n}^{p}$. The convergence of $S_{2}$ could then be proved in the same way.

This line of argument leads up to and is completed by the following theorem.
326. If $p>1, a_{n} \geqq 0$, and $A_{n}=a_{1}+a_{2}+\ldots+a_{n}$, then

$$
\begin{equation*}
\Sigma\left(\frac{A_{n}}{n}\right)^{p}<\left(\frac{p}{p-1}\right)^{p} \Sigma a_{n}^{p} \tag{9.8.1}
\end{equation*}
$$

unless all the a are zero. The constant is the best possible.

[^61]The corresponding theorem for integrals is
327.a If $p>1, f(x) \geqq 0$, and $F(x)=\int_{0}^{x} f(t) d t$, then

$$
\begin{equation*}
\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x<\left(\frac{p}{p-1}\right)^{p} \int_{0}^{\infty} f^{p} d x \tag{9.8.2}
\end{equation*}
$$

unless $f \equiv 0$. The constant is the best possible.
These theorems were first proved by Hardy (2), except that Hardy was unable to fix the constant in Theorem 326. This imperfection was removed by Landau (4). A great many alternative proofs of the theorems have been given by various writers, for example by Broadbent (1), Elliott (1), Grandjot (1), Hardy (4), Kaluza and Szegö (1), Knopp (1). We begin by giving Elliott's proof of Theorem 326 and Hardy's proof of Theorem 327. ${ }^{\text {b }}$
(i) In proving Theorem 326 we may suppose that $a_{1}>0$. For if we suppose that $a_{1}=0$, and replace $a_{n+1}$ by $b_{n},(9.8 .1)$ becomes

$$
\left(\frac{b_{1}}{2}\right)^{p}+\left(\frac{b_{1}+b_{2}}{3}\right)^{p}+\ldots<\left(\frac{p}{p-1}\right)^{p}\left(b_{1}^{p}+b_{2}^{p}+\ldots\right)
$$

an inequality weaker than (9.8.1) itself.
Let us write $\alpha_{n}$ for $A_{n} / n$, and agree that any number with suffix 0 is 0 . We have then

$$
\begin{aligned}
& \alpha_{n}^{p}- \frac{p}{p-1} \alpha_{n}^{p-1} a_{n}=\alpha_{n}^{p}-\frac{p}{p-1}\left\{n \alpha_{n}-(n-1) \alpha_{n-1}\right\} \alpha_{n}^{p-1} \\
&=\alpha_{n}^{p}\left(1-\frac{n p}{p-1}\right)+\frac{(n-1) p}{p-1} \alpha_{n}^{p-1} \alpha_{n-1} \\
& \leqq \alpha_{n}^{p}\left(1-\frac{n p}{p-1}\right)+\frac{n-1}{p-1}\left\{(p-1) \alpha_{n}^{p}+\alpha_{n-1}^{p}\right\}^{\mathrm{c}} \\
&= \frac{1}{p-1}\left\{(n-1) \alpha_{n-1}^{p}-n \alpha_{n}^{p}\right\} . \\
& \text { ace } \quad \sum_{1}^{N} \alpha_{n}{ }^{p}-\frac{p}{p-1} \sum_{1}^{N} \alpha_{n}^{p-1} a_{n} \leqq-\frac{N \alpha_{N}^{p}}{p-1} \leqq 0 ;
\end{aligned}
$$

Hence
a We have already encountered this theorem in Ch. VII, but the proof which we gave there (in detail only when $p=2$ ) was intended primarily as an illustration of variational methods and has no particular pretensions to simplicity.
${ }^{b}$ We have expanded the proofs so as to deal with the question of equality.
c By Theorem 9.
and therefore, by Theorem 13,

$$
\begin{equation*}
\sum_{1}^{N} \alpha_{n}{ }^{p} \leqq \frac{p}{p-1} \sum_{1}^{N} \alpha_{n}^{p-1} a_{n} \leqq \frac{p}{p-1}\left(\sum_{1}^{N} a_{n}^{p}\right)^{1 / p}\left(\sum_{1}^{N} \alpha_{n}^{p}\right)^{1 / p^{\prime}} \tag{9.8.3}
\end{equation*}
$$

Dividing by the last factor on the right (which is certainly positive), and raising the result to the $p$ th power, we obtain

$$
\begin{equation*}
\sum_{1}^{N} \alpha_{n}{ }^{p} \leqq\left(\frac{p}{p-1}\right)^{p} \sum_{1}^{N} a_{n}{ }^{p} . \tag{9.8.4}
\end{equation*}
$$

When we make $N$ tend to infinity we obtain (9.8.1), except that we have ' $\leqq$ ' in place of ' $<$ '. In particular we see that $\Sigma \alpha_{n}{ }^{p}$ is finite.

Returning to (9.8.3), and replacing $N$ by $\infty$, we obtain

$$
\begin{equation*}
\Sigma \alpha_{n}^{p} \leqq \frac{p}{p-1} \Sigma \alpha_{n}^{p-1} a_{n} \leqq \frac{p}{p-1}\left(\Sigma a_{n}^{p}\right)^{1 / p}\left(\Sigma \alpha_{n}^{p}\right)^{1 / p^{\prime}} \tag{9.8.5}
\end{equation*}
$$

There is inequality in the second place unless $\left(a_{n}{ }^{p}\right)$ and ( $\alpha_{n}{ }^{p}$ ) are proportional, i.e. unless $a_{n}=C \alpha_{n}$, where $C$ is independent of $n$. If this is so then (since $a_{1}=\alpha_{1}>0$ ) $C$ must be 1 , and then - $A_{n}=n a_{n}$ for all $n$. This is only possible if all the $a$ are equal, and this is inconsistent with the convergence of $\Sigma a_{n}{ }^{p}$. Hence

$$
\begin{equation*}
\Sigma \alpha_{n}^{p}<\frac{p}{p-1}\left(\Sigma{a_{n}}^{p}\right)^{1 / p}\left(\Sigma \alpha_{n}^{p}\right)^{1 / p^{\prime}} \tag{9.8.6}
\end{equation*}
$$

and (9.8.1) follows from (9.8.6) as (9.8.4) followed from (9.8.3).
To prove the constant factor the best possible, we take

Then

$$
\begin{gathered}
a_{n}=n^{-1 / p}(n \leqq N), \quad a_{n}=0(n>N) \\
\Sigma a_{n}^{p}=\sum_{1}^{N} \frac{1}{n} \\
A_{n}=\sum_{1}^{n} \nu^{-1 / p}>\int_{1}^{n} x^{-1 / p} d x=\frac{p}{p-1}\left\{n^{(p-1) / p}-1\right\} \quad(n \leqq N) \\
\left(\frac{A_{n}}{n}\right)^{p}>\left(\frac{p}{p-1}\right)^{p} \frac{1-\epsilon_{n}}{n} \quad(n \leqq N)
\end{gathered}
$$

where $\epsilon_{n} \rightarrow 0$ when $n \rightarrow \infty$. It follows that

$$
\Sigma\left(\frac{A_{n}}{n}\right)^{p}>\sum_{1}^{N}\left(\frac{A_{n}}{n}\right)^{p}>\left(\frac{p}{p-1}\right)^{p}\left(1-\eta_{N}\right) \Sigma a_{n}^{p}
$$

where $\eta_{N} \rightarrow 0$ when $N \rightarrow \infty$. Hence any inequality of the type

$$
\Sigma\left(\frac{A_{n}}{n}\right)^{p}<\left(\frac{p}{p-1}\right)^{p}(1-\epsilon) \Sigma a_{n}^{p}
$$

is false if $a_{n}$ is chosen as above and $N$ is sufficiently large.
An alternative procedure is to take $a_{n}=n^{-(1 / p)-\epsilon}$ for all $n$, and to make $\epsilon$ small. Compare $\S 9.5$ : it was this procedure which we followed there.
(ii) If $0<\xi<X$, we have

$$
\begin{aligned}
& \int_{\xi}^{X}\left(\frac{F}{x}\right)^{p} d x=-\frac{1}{p-1} \int_{\xi}^{X} F^{p} \frac{d}{d x} x^{1-p} d x \\
& \quad=\frac{\xi^{1-p} F^{p}(\xi)}{p-1}-\frac{X^{1-p} F^{p}(X)}{p-1}+\frac{p}{p-1} \int_{\xi}^{X} x^{1-p} F^{p-1} f d x .
\end{aligned}
$$

But, by Theorem 222, $\xi^{1-p} F^{p}(\xi) \rightarrow 0$ when $f^{p}$ is integrable and $\xi \rightarrow 0$. Hence

$$
\begin{align*}
\int_{0}^{X}\left(\frac{F}{x}\right)^{p} d x & \leqq \frac{p}{p-1} \int_{0}^{x}\left(\frac{F}{x}\right)^{p-1} f d x  \tag{9.8.7}\\
& \leqq \frac{p}{p-1}\left\{\int_{0}^{X}\left(\frac{F}{x}\right)^{p} d x\right\}^{1 / p^{\prime}}\left(\int_{0}^{X} f^{p} d x\right)^{1 / p}
\end{align*}
$$

If $f$ is not nul in $(0, X)$, the left-hand side of (9.8.7) is positive. Hence (9.8.7) gives

$$
\begin{equation*}
\int_{0}^{X}\left(\frac{\bar{F}}{x}\right)^{p} d x \leqq\left(\frac{p}{p-1}\right)^{p} \int_{0}^{X} f^{p} d x \tag{9.8.8}
\end{equation*}
$$

and when we make $X \rightarrow \infty$ we obtain (9.8.2), except that ' $<$ ' is replaced by ' $\leqq$ '. In particular, the integral on the left-hand side of (9.8.2) is finite.

It follows that all the integrals in (9.8.7) remain finite when $X$ is replaced by $\infty$, and that

$$
\begin{align*}
\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x & \leqq \frac{p}{p-1} \int_{0}^{\infty}\left(\frac{F}{x}\right)^{p-1} f d x  \tag{9.8.9}\\
& \leqq \frac{p}{p-1}\left\{\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x\right\}^{1 / p^{\prime}}\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}
\end{align*}
$$

The last sign of inequality may be replaced by ' $<$ ' unless $x^{-p} F^{p}$

[^62]and $f^{p}$ are effectively proportional. This would make $f$ a power of $x$, and then $\int f^{p} d x$ would be divergent.

Hence

$$
\begin{equation*}
\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x<\frac{p}{p-1}\left\{\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x\right\}^{1 / p^{\prime}}\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p} \tag{9.8.10}
\end{equation*}
$$

unless $f$ is nul. Since the integral on the left-hand side is positive and finite, (9.8.2) now follows from (9.8.10) as (9.8.8) followed from (9.8.7).

The proof that the constant is the best possible follows the same lines as before: take $f(x)=0$ for $x<1, f(x)=x^{-(1 / p)-\epsilon}$ for $x \geqq 1$.

Elliott's proof of Theorem 326 applies to Theorem 327 also, with the obvious modifications. The proof of Theorem 327 given in (ii) may be adapted to series, but does not give the best possible value of the constant.
(iii) The following proof of Theorem 327 (due to Ingham) is also interesting: we shall be content with proving the form with ' $\leqq$ '. We use Theorem 203, supposing that the intervals of integration are each ( 0,1 ), that the weight-functions are 1, and that

Then

$$
r=1, \quad s=p>1, \quad f(x, y)=f(x y)
$$

$$
\mathfrak{M}_{1}^{(x)} f(x y)=\int_{0}^{1} f(x y) d x=\frac{F(y)}{y}
$$

$$
\mathfrak{M}_{p}^{(y)} f(x y)=\left\{\int_{0}^{1} f^{p}(x y) d y\right\}^{1 / p}=\left\{\frac{1}{x} \int_{0}^{x} f^{p}(t) d t\right\}^{1 / p} \leqq\left\{\frac{1}{x} \int_{0}^{1} f^{p}(t) d t\right\}^{1 / p}
$$

for $x \leqq 1$. Hence, by Theorem 203,

$$
\left\{\int_{0}^{1}\left(\frac{F}{y}\right)^{p} d y\right\}^{1 / p} \leqq\left(\int_{0}^{1} f^{p} d t\right)^{1 / p} \int_{0}^{1} x^{-1 / p} d x=\frac{p}{p-1}\left(\int_{0}^{1} f^{p} d x\right)^{1 / p}
$$

We then obtain the result by putting

$$
x=X / c, \quad f(X / c)=g(X),
$$

replacing $X, g$ by $x, f$, and making $c \rightarrow \infty$.
9.9. Further integral inequalities. There are many analogues and extensions of Theorems 326 and 327, which have been proved by different writers in different ways; and we give some of these theorems here. We consider integral inequalities first, since we can derive most of these in a simple and uniform manner from Theorem 319, and the corresponding theorems for series sometimes involve slight additional complications.
(1) Take, in Theorem 319,

$$
K(x, y)=1 / y \quad(x \leqq y), \quad K(x, y)=0(x>y)
$$

Then, if $p>1$,

$$
k=\int_{0}^{\infty} K(x, 1) x^{-1 / p} d x=\int_{0}^{1} x^{-1 / p} d x=\frac{p}{p-1}
$$

and all the conditions on $K$ are satisfied. Hence (b) and (c) of Theorem 319 give

$$
\begin{equation*}
\int_{0}^{\infty} d y\left(\frac{1}{y} \int_{0}^{y} f(x) d x\right)^{p} \leqq\left(\frac{p}{p-1}\right)^{p} \int_{0}^{\infty} f^{p} d x \tag{9.9.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{0}^{\infty} d x\left(\int_{x}^{\infty} \frac{g(y)}{y} d y\right)^{p^{\prime}} \leqq\left(\frac{p}{p-1}\right)^{\boldsymbol{p}^{\prime}} \int_{0}^{\infty} g^{p^{\prime}} d y \tag{9.9.2}
\end{equation*}
$$

Of these inequalities, (9.9.1) is Theorem 327, with ' $\leqq$ ' for ' <'; we cannot quote ' $<$ ' from the general theorem because $K$ is not always positive. But equality in (9.9.1), with a non-nul $f$, involves

$$
\iint K(x, y) f(x) g(y) d x d y=\frac{p}{p-1}\left(\int f^{p} d x\right)^{1 / p}\left(\int g^{p^{\prime}} d x\right)^{1 / p^{\prime}}
$$

with non-nul $f$ and $g$. The argument of $\S 9.3$ then shows that (9.3.1) is true for $x<y$, and that $f \equiv C x^{-1 / p}$ for small $x$, which is inconsistent with the convergence of $\int f^{p} d x$.

Similarly we can prove that there is inequality in (9.9.2) unless $g$ is nul. A trivial transformation then gives
328. If $p>1$ and

$$
F(x)=\int_{x}^{\infty} f(t) d t
$$

then

$$
\begin{equation*}
\int_{0}^{\infty} F^{p} d x<p^{p} \int_{0}^{\infty}(x f)^{p} d x \tag{9.9.3}
\end{equation*}
$$

unless $f \equiv 0$. The constant is the best possible.
(2) More generally, take

$$
K(x, y)=\frac{1}{\Gamma(r)} \frac{(y-x)^{r-1}}{y^{r}}(x<y), \quad K(x, y)=0 \quad(x \geqq y)
$$

with $r>0$. With $r=1$, we come back to (1). We now have

$$
k=\frac{1}{\Gamma(r)} \int_{0}^{1} x^{-1 / p}(1-x)^{r-1} d x=\frac{\Gamma\left(1-\frac{1}{p}\right)}{\Gamma\left(r+1-\frac{1}{p}\right)}
$$

We are thus led to
329. If $p>1, r>0$, and

$$
\begin{equation*}
f_{r}(x)=\frac{1}{\Gamma(r)} \int_{0}^{x}(x-t)^{r-1} f(t) d t, \tag{9.9.4}
\end{equation*}
$$

then

$$
\begin{equation*}
\int_{0}^{\infty}\left(\frac{f_{r}}{x^{r}}\right)^{p} d x<\left\{\frac{\Gamma\left(1-\frac{1}{p}\right)}{\Gamma\left(r+1-\frac{1}{p}\right)}\right\}^{p} \int_{0}^{\infty} f^{p} d x \tag{9.9.5}
\end{equation*}
$$

unless $f \equiv 0$. If

$$
\begin{equation*}
f_{r}(x)=\frac{1}{\Gamma(r)} \int_{x}^{\infty}(t-x)^{r-1} f(t) d t, \tag{9.9.6}
\end{equation*}
$$

then

$$
\begin{equation*}
\int_{0}^{\infty} f_{r}^{p} d x<\left\{\frac{\Gamma\left(\frac{1}{p}\right)}{\Gamma\left(r+\frac{1}{p}\right)}\right\}^{p} \int_{0}^{\infty}\left(x^{r} f\right)^{p} d x \tag{9.9.7}
\end{equation*}
$$

unless $f \equiv 0$. In each case the constant is the best possible.
The function $f_{r}(x)$ of (9.9.4) is the 'Riemann-Liouvilleintegral'a of $f(x)$ of order $r$, with 'origin' 0 . The function (9.9.6) is the 'Weyl integral' of order $r$, which is in some ways more convenient, especially in the theory of Fourier series.
(3) Take

$$
K(x, y)=\frac{y^{a-1}}{x^{a}}(x \leqq y), \quad K(x, y)=0 \quad(x>y),
$$

with $\alpha<1 / p^{\prime}$. Then

$$
k=\int_{0}^{1} x^{-a-1 / p} d x=\frac{p}{p-p \alpha-1},
$$

and (b) and (c) of Theorem 319 give
(9.9.8) $\int_{0}^{\infty} y^{p(a-1)}\left(\int_{0}^{y} x^{-\alpha} f(x) d x\right)^{p} d y<\left(\frac{p}{p-\alpha p-1}\right)^{p} \int_{0}^{\infty} f^{p} d x$,
(9.9.9) $\int_{0}^{\infty} x^{-\alpha p^{\prime}}\left(\int_{x}^{\infty} y^{\alpha-1} g(y) d y\right)^{p \prime} d x<\left(\frac{p^{\prime}}{1-\alpha p^{\prime}}\right)^{p^{\prime}} \int_{0}^{\infty} g^{p^{\prime}} d y$.

Changing the notation, we obtain
330. If $p>1, r \neq 1$, and $F(x)$ is defined by

$$
F(x)=\int_{0}^{x} f(t) d t(r>1), \quad F(x)=\int_{x}^{\infty} f(t) d t(r<1),
$$

[^63]then
\[

$$
\begin{equation*}
\int_{0}^{\infty} x^{-r} F^{p} d x<\left(\frac{p}{|r-1|}\right)^{p} \int_{0}^{\infty} x^{-r}(x f)^{p} d x \tag{9.9.10}
\end{equation*}
$$

\]

unless $f \equiv 0 .{ }^{\text {a }}$
The constant is the best possible. It is also easy to verify that, when $p=1$, the two sides of (9.9.10) are equal.
9.10. Further theorems concerning series. Among the analogues and extensions of Theorem 326 we select the following.
(1) The theorem related to Theorem $\mathbf{3 2 6}$ as Theorem $\mathbf{3 2 8}$ is to Theorem 327 is
331. ${ }^{\text {b }}$ If $p>1$, then

$$
\Sigma\left(a_{n}+a_{n+1}+\ldots\right)^{p}<p^{p} \Sigma\left(n a_{n}\right)^{p},
$$

unless $\left(a_{n}\right)$ is nul. The constant is the best possible.
This theorem is 'reciprocal' to Theorem 326 in the sense of $\S 9.7$ (2), i.e. deducible from the latter theorem by the converse of Hölder's inequality. It may be instructive to set out the proof in detail, although what we say amounts to a repetition in a special case of what we have explained more generally before ${ }^{c}$.

If $K(x, y)$ is defined as in $\S 9.9(1)$, then

$$
\begin{align*}
& \Sigma \Sigma K(m, n) a_{m} b_{n}=\Sigma \sum_{m \leqq n} \frac{a_{m} b_{n}}{n}=\Sigma \frac{a_{1}+a_{2}+\ldots+a_{n}}{n} b_{n}=\Sigma \frac{A_{n}}{n} b_{n}  \tag{9.10.1}\\
& \leqq\left\{\Sigma\left(\frac{A_{n}}{n}\right)^{p}\right\}^{1 / p}\left(\Sigma b_{n}{ }^{p^{\prime}}\right)^{1 / p^{\prime}}<\frac{p}{p-1}\left(\Sigma a_{m}^{p}\right)^{1 / p}\left(\Sigma b_{n}^{p^{\prime}}\right)^{1 / p^{\prime}}
\end{align*}
$$

by Theorems 13 and 326, unless (a) or (b) is nul.
On the other hand

$$
\Sigma \Sigma K(m, n) a_{m} b_{n}=\Sigma a_{m}\left(\frac{b_{m}}{m}+\frac{b_{m+1}}{m+1}+\ldots\right) ;
$$

and the maximum of this, for all (a) for which $\Sigma a_{m}{ }^{p}=1$, is, by Theorem 15,

$$
\left\{\Sigma\left(\frac{b_{m}}{m}+\frac{b_{m+1}}{m+1}+\ldots\right)^{p^{\prime}}\right\}^{1 / p^{\prime}}
$$

[^64]Hence, by (9.10.1),

$$
\left\{\Sigma\left(\frac{b_{m}}{m}+\frac{b_{m+1}}{m+1}+\ldots\right)^{p^{\prime}}\right\}^{1 / p^{\prime}}<\frac{p}{p-1}\left(\Sigma b_{n} p^{\prime}\right)^{1 / p^{\prime}}=p^{\prime}\left(\Sigma b_{n} p^{\prime}\right)^{1 / p^{\prime}}
$$

Changing $b_{m}$ into $m a_{m}$, and $p^{\prime}$ into $p$, we obtain Theorem 331.
That the constant $p^{p}$ is the best possible follows from the last remark of § 9.7 (2).
(2) 332. If $p>1, a_{n} \geqq 0, \lambda_{n}>0$, and

$$
\Lambda_{n}=\lambda_{1}+\lambda_{2}+\ldots+\lambda_{n}, \quad A_{n}=\lambda_{1} a_{1}+\lambda_{2} a_{2}+\ldots+\lambda_{n} a_{n}
$$

then

$$
\Sigma \lambda_{n}\left(\frac{A_{n}}{\Lambda_{n}}\right)^{p}<\left(\frac{p}{p-1}\right)^{p} \Sigma \lambda_{n} a_{n}^{p}
$$

unless $\left(a_{n}\right)$ is nul.
This theorem, which is related to Theorem 326 as Theorem 321 is related to Theorem 315, may be proved in various ways. In the first place, itmay be deduced from Theorem 320 by a specialisation of $K$ (as Theorem 327 was deduced from Theorem 319 in §9.9); but the question of possible equality then needs a little attention. Perhaps the simplest proof is by a direct adaptation of Elliott's argument in §9.8. If $\alpha_{n}=A_{n} / \Lambda_{n}$, we find that

$$
\lambda_{n} \alpha_{n}^{p}-\frac{p}{p-1} \lambda_{n} \alpha_{n}^{p-1} a_{n} \leqq \frac{1}{p-1}\left(\Lambda_{n-1} \alpha_{n-1}^{p}-\Lambda_{n} \alpha_{n}^{p}\right)
$$

and the proof may be completed as in §9.8. ${ }^{\text {a }}$
The theorem may also be deduced from Theorem 327 by taking $f$ to be an appropriate step-function (the process by which, in §9.4, we deduced Theorem 320 from Theorem 319). We shall not set this out in detail ${ }^{\text {b }}$; but the remark raises questions of which we say something more in the next section.

### 9.11. Deduction of theorems on series from theorems

 on integrals ${ }^{c}$. The process of deduction just referred to, and actually used in $\S 9.4$, is very natural and often effective. It is however apt to lead to difficulties of detail, so that direct methods are usually preferable. We illustrate this by giving a deduction of[^65]Theorem 326 from Theorem 327, which leads us incidentally to a remark of considerable intrinsic interest.

We observe first that it is sufficient to prove Theorem 326 on the hypothesis that $a_{n}$ decreases as $n$ increases. This follows from a theorem which is of sufficient interest to be stated separately.
333. If the $a_{n}$ are given except in arrangement, and $\phi(u)$ is a positive increasing function of $u$, then

$$
\Sigma \phi\left(\frac{A_{n}}{n}\right)
$$

is greatest when the $a_{n}$ are arranged in decreasing order.
To prove Theorem 333 we have only to observe that, if $\nu>\mu$ and $a_{\nu}>a_{\mu}$, the effect of exchanging $a_{\mu}$ and $a_{\nu}$ is to leave $A_{n}$ unchanged when $n<\mu$ or $n \geqq \nu$, and to increase $A_{n}$ when $\mu \leqq n<\nu$. The theorem is one of a type which we shall discuss in much greater detail in Ch. X .

Suppose now that $a_{n}$ decreases, and that Theorem 327 has been proved. We define $f(x)$ by

$$
f(x)=a_{n} \quad(n-1 \leqq x<n) .
$$

Then

$$
\begin{equation*}
\Sigma a_{n}^{p}=\int_{0}^{\infty} f^{p} d x \tag{9.11.1}
\end{equation*}
$$

If $n<x<n+1$, then

$$
\frac{F(x)}{x}=\frac{a_{1}+a_{2}+\ldots+a_{n}+(x-n) a_{n+1}}{x}=\frac{A_{n}-n a_{n+1}+x a_{n+1}}{x}
$$

and

$$
A_{n}-n a_{n+1} \geqq 0
$$

so that $F / x$ decreases from $A_{n} / n$ to $A_{n+1} /(n+1)$ when $x$ increases from $n$ to $n+1$. Hence

$$
\frac{F}{x} \geqq \frac{A_{n+1}}{n+1} \quad(n<x<n+1)
$$

and so

$$
\begin{equation*}
\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x \geqq \sum_{1}^{\infty}\left(\frac{A_{n}}{n}\right)^{p} \tag{9.11.2}
\end{equation*}
$$

Theorem 326 now follows from (9.11.1); (9.11.2), and Theorem 333.

If the reader will try to deduce Theorem 331 from Theorem 328 similarly, he will find some difficulty. Something is lost in the passage from integrals to series, and it is by no means always that (as here) the passage can be made without damage to the final result.
9.12. Carleman's inequality. If in Theorem 326 we write $a_{n}$ for $a_{n}{ }^{p}$, we obtain

$$
\begin{equation*}
\text { 1) } \Sigma\left(\frac{a_{1}^{1 / p}+a_{2}{ }^{1 / p}+\ldots+a_{n}^{1 / p}}{n}\right)^{p}<\left(\frac{p}{p-1}\right)^{p} \Sigma a_{n} \text {. } \tag{9.12.1}
\end{equation*}
$$

If we make $p \rightarrow \infty$, and use Theorem 3, we obtain

$$
\Sigma\left(a_{1} a_{2} \ldots a_{n}\right)^{1 / n} \leqq e \Sigma a_{n}
$$

and this suggests the more complete theorem which follows.
334.a

$$
\Sigma\left(a_{1} a_{2} \ldots a_{n}\right)^{1 / n}<e \Sigma a_{n},
$$

unless $\left(a_{n}\right)$ is nul. The constant is the best possible.
It is natural to attempt to prove the complete theorem by means of Theorem 9; but a direct application of Theorem 9 to the left-hand side of (9.12.1) is insufficient ${ }^{b}$. To remedy this, we apply Theorem 9 not to $a_{1}, a_{2}, \ldots, a_{n}$ but to $c_{1} a_{1}, c_{2} a_{2}, \ldots, c_{n} a_{n}$, and choose the $c$ so that, when $\Sigma a_{n}$ is near the boundary of convergence, these numbersshall be 'roughly equal'. This requires that $c_{n}$ shall be roughly of the order of $n$.
These considerations suggest the following proof. We have

$$
\begin{gathered}
\Sigma\left(a_{1} a_{2} \ldots a_{n}\right)^{1 / n}=\Sigma\left(\frac{c_{1} a_{1} \cdot c_{2} a_{2} \ldots c_{n} a_{n}}{c_{1} c_{2} \ldots c_{n}}\right)^{1 / n} \\
\leqq \sum_{n}\left(c_{1} c_{2} \ldots c_{n}\right)^{-1 / n} \frac{1}{n} \sum_{m \leqq n} c_{m} a_{m} \\
=\sum_{m} a_{m} c_{m} \sum_{n \geqq m} \frac{1}{n}\left(c_{1} c_{2} \ldots c_{n}\right)^{-1 / n}
\end{gathered}
$$

[^66]In order that the inner summation should be easily effected, we choose

$$
\left(c_{1} c_{2} \ldots c_{n}\right)^{1 / n}=n+1
$$

when

$$
c_{m}=\frac{(m+1)^{m}}{m^{m-1}}, \quad \sum_{n \geqq m} \frac{1}{n}\left(c_{1} c_{2} \ldots c_{n}\right)^{-1 / n}=\sum_{n \geqq m} \frac{1}{n(n+1)}=\frac{1}{m},
$$

and then

$$
\Sigma\left(a_{1} a_{2} \ldots a_{n}\right)^{1 / n} \leqq \Sigma \frac{a_{m} c_{m}}{m}=\Sigma a_{m}\left(1+\frac{1}{m}\right)^{m}<e \Sigma a_{m},
$$

by Theorem 140 , unless $a_{m}$ is nul.
We can prove the constant best possible as in § 9.5. We may, for example, take $a_{n}=1 / n$ for $n \leqq \mu, a_{n}=0$ for $n>\mu$, and make $\mu$ tend to infinity.

The corresponding integral theorem is
335.a If $f$ is not nul, then

$$
\int_{0}^{\infty} \exp \left\{\frac{1}{x} \int_{0}^{x} \log f(t) d t\right\} d x<e \int_{0}^{\infty} f(x) d x
$$

9.13. Theorems with $0<p<1$. We have supposed so far that the parameter $p$ involved in our theorems is greater than 1. A good many of them, however, have analogues with a $p$ less than 1 , and we give a selection of them in this section. The characteristic difference between the two cases lies (as is to be expected after our experience with Hölder's and Minkowski's inequalities) in a reversal of the sign of inequality.
(1) 336. If $K(x, y)$ is non-negative and homogeneous of degree $-1,0<p<1$, and

$$
\int_{0}^{\infty} K(x, 1) x^{-1 / p} d x=\int_{0}^{\infty} K(1, y) y^{-1 / p^{\prime}} d y=k<\infty
$$

then
(a) $\int_{0}^{\infty} \int_{0}^{\infty} K(x, y) f(x) g(y) d x d y \geqq k\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{p^{\prime}} d y\right)^{1 / p^{\prime}}$,

$$
\begin{equation*}
\int_{0}^{\infty} d y\left(\int_{0}^{\infty} K(x, y) f(x) d x\right)^{p} \geqq k^{p} \int_{0}^{\infty} f^{p}(x) d x \tag{b}
\end{equation*}
$$

Here, in accordance with the conventions of $\S \S 5.1$ and 6.5 ,
(a) means 'if the double integral, and the second integral on the right, are finite, then the first integral on the right is finite, and ...'; and (b) means 'if the integral on the left is finite, then the integral on the right is finite, and ...'.

If we use the second method of $\S 9.3$, the proof of $(a)$ is the same as that of $(a)$ of Theorem 319. The sign is reversed because we use Hölder's inequality with $p<1$. To deduce (b) from (a), we appeal to Theorem 234. We leave it to the reader to frame the corresponding theorem for $p<0$, and to consider the question of equality.

We cannot take $K=1 /(x+y)$, since then $k=\infty$. There is therefore no exact analogue of Hilbert's theorem.
(2) 337. If $0<p<1, f(x) \geqq 0$,
and

$$
\int_{0}^{\infty} f^{p} d x<\infty
$$

$$
F(x)=\int_{x}^{\infty} f(t) d t
$$

then

$$
\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x>\left(\frac{p}{1-p}\right)^{p} \int_{0}^{\infty} f^{p} d x
$$

unless $f \equiv 0$. The constant is the best possible.
We may deduce Theorem 337, in an imperfect form, from Theorem 336, by taking
when

$$
\begin{gathered}
K(x, y)=0 \quad(x<y), \quad K(x, y)=\frac{1}{y}(x \geqq y) \\
k=\int_{1}^{\infty} x^{-1 / p} d x=\frac{p}{1-p}
\end{gathered}
$$

To prove the complete theorem in this way would involve a discussion of the sign of inequality in Theorem 336 (and so in Theorem 234). We therefore follow a direct method analogous to that of §9.8.

We may suppose

$$
\int_{x}^{\infty} f(t) d t, \quad \int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x
$$

finite, since otherwise there would be nothing to prove.

We have
(9.13.1)

$$
\int_{\xi}^{X}\left(\frac{\boldsymbol{F}}{x}\right)^{p} d x=\frac{1}{1-p}\left[x^{1-p} F^{p}(x)\right]_{\xi}^{X}+\frac{p}{1-p} \int_{\xi}^{X}\left(\frac{F}{x}\right)^{p-1} f d x
$$

Since $F$ decreases as $x$ increases,

$$
x^{1-p} F^{p}(x)=2\left(\frac{F(x)}{x}\right)^{p} \frac{1}{2} x \leqq 2 \int_{\frac{1}{\frac{1}{2}}}^{x}\left(\frac{F}{t}\right)^{p} d t
$$

tends to 0 both when $x \rightarrow 0$ and when $x \rightarrow \infty$. Hence (9.13.1) gives in the limit

$$
\int_{0}^{\infty}\left(\frac{F}{x}\right)^{p} d x=\frac{p}{1-p} \int_{0}^{\infty}\left(\frac{F}{x}\right)^{p-1} f d x
$$

and the proof may be completed as before.
For a more complete result, corresponding to Theorem 330, see Theorem 347.
(3) Finally we prove a theorem which is related to Theorem 326 roughly as Theorem 337 is related to Theorem327. The correspondence is not quite precise, and the theorem illustrates very instructively the slight additional complications which are sometimes inherent in a theorem concerning series.
338. ${ }^{\text {a }}$ If $0<p<1$ and $\Sigma a_{n}{ }^{p}<\infty$, then

$$
\Sigma^{\prime}\left(\frac{a_{n}+a_{n+1}+\ldots}{n}\right)^{p}>\left(\frac{p}{1-p}\right)^{p} \Sigma a_{n}^{p},
$$

unless $\left(a_{n}\right)$ is nul. The dash over the summation on the left-hand side implies that the term for which $n=1$ is to be multiplied by

$$
1+\frac{1}{1-p}
$$

The constant is the best possible.
In Theorem 337, take

$$
f(x)=0(0<x<1), \quad f(x)=a_{n}(0<n \leqq x<n+1) .
$$

Then, if $0<n \leqq x<n+1$,

$$
\frac{F}{x}=\frac{(n+1-x) a_{n}+a_{n+1}+\ldots}{x} \leqq \frac{a_{n}+a_{n+1}+\ldots}{n} .
$$

Hence

$$
\int_{1}^{\infty}\left(\frac{F}{x}\right)^{p} d x \leqq \sum_{1}^{\infty}\left(\frac{a_{n}+a_{n+1}+\ldots}{n}\right)^{p},
$$

while $\quad \int_{0}^{1}\left(\frac{F}{x}\right)^{p} d x=\left(a_{1}+a_{2}+\ldots\right)^{p} \int_{0}^{1} x^{-p} d x=\frac{\left(a_{1}+a_{2}+\ldots\right)^{p}}{1-p}$;
and the result follows from Theorem 337.
${ }^{\text {a }}$ The substance of this theorem was communicated to us by Prof. Elliott in 1927.

Some such gloss as that contained in the last clause of the theorem is necessary; the result is not necessarily true if the dash is omitted ${ }^{\text {a }}$.
9.14. A theorem with two parameters $p$ and $q$. We conclude this chapter with a theorem which, although again an extension of Hilbert's theorem, has peculiarities which do not occur in any of the earlier theorems of the chapter. It involves two independent indices $p$ and $q$ and an undetermined constant $K(p, q)$.
339. If

$$
p>1, \quad q>1, \quad \frac{1}{p}+\frac{1}{q} \geqq 1,
$$

so that

$$
0<\lambda=2-\frac{1}{p}-\frac{1}{q}=\frac{1}{p^{\prime}}+\frac{1}{q^{\prime}} \leqq 1,
$$

$$
\sum_{1}^{\infty} \sum_{1}^{\infty} \frac{a_{m} b_{n}}{(m+n)^{\lambda}} \leqq K\left(\sum_{1}^{\infty} a_{m}{ }^{p}\right)^{1 / p}\left(\sum_{1}^{\infty} b_{n}{ }^{q}\right)^{1 / q},
$$

where $K=K(p, q)$ depends on $p$ and $q$ only.
This theorem reduces to Theorem 315 when $q=p^{\prime}, \lambda=1$ : in that case we know the best possible value of $K$. The best value has not been found in the general case, and the problem of determining it appears to be difficult. We shall prove later (§10.17) a deeper theorem in which $\lambda<1$ and $m+n$ is replaced by $|m-n|$ (equal values being then excluded from the summation).

It is sufficient to prove that, if $\Sigma a_{m}{ }^{p}=A, \Sigma b_{n}{ }^{q}=B$, then

$$
\begin{equation*}
\sum_{m} a_{m} \sum_{n \leqq m} \frac{b_{n}}{(m+n)^{\lambda}} \leqq K A^{1 / p} B^{1 / q}, \tag{9.14.1}
\end{equation*}
$$

and for this, by Theorem 13, that

$$
\begin{gather*}
{ }_{m} \beta_{m}{ }^{p^{\prime} \leqq K B^{p^{\prime} q},}  \tag{9.14.2}\\
\beta_{m}=\sum_{n \leqq m} \frac{b_{n}}{(m+n)^{\lambda}} .
\end{gather*}
$$

where
Now

$$
\beta_{m} \leqq m^{-\lambda} \sum_{n \leqq m} b_{n}=m^{-\lambda} B_{m}
$$

and $p^{\prime} \geqq q$. Hence

$$
\Sigma \beta_{m}{ }^{p^{\prime}} \leqq \Sigma m^{-p^{\prime} \lambda} B_{m}{ }^{p^{\prime}}=\Sigma\left(\frac{B_{m}}{m}\right)^{q} B_{m}^{p^{\prime}-q} m^{q-p^{\prime} \lambda} .
$$

[^67]But

$$
B_{m}=\sum_{1}^{m} b_{n} \leqq m^{1 / q^{\prime}}\left(\sum_{1}^{m} b_{n}^{q}\right)^{1 / q} \leqq B^{1 / q} m^{1 / q^{\prime}}
$$

and

$$
\frac{p^{\prime}-q}{q^{\prime}}+q-p^{\prime} \lambda=0
$$

Hence

$$
\Sigma \beta_{m}^{p^{\prime}} \leqq B^{\left(p^{\prime}-q\right) / q} \Sigma\left(\frac{B_{m}}{m}\right)^{q} \leqq\left(\frac{q}{q-1}\right)^{q} B^{\left(p^{\prime}-q\right) / q+1}=K B^{p^{\prime} / q}
$$

by Theorem 326. This proves (9.14.2).
Similarly we prove
340. Under the same conditions as those of Theorem 339

$$
\int_{0}^{\infty} \int_{0}^{\infty} \frac{f(x) g(y)}{(x+y)^{\lambda}} d x d y \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{q} d y\right)^{1 / q}
$$

## MISCELLANEOUS THEOREMS AND EXAMPLES

341. If (i) $a_{m}, b_{n}, f(x), g(y)$ are non-negative, (ii) the summations go from 1 to $\infty$ and the integrations from 0 to $\infty$,
(iii) $\quad\left(\Sigma a_{m}{ }^{p}\right)^{1 / p}=A, \quad\left(\Sigma b_{n}{ }^{p \prime}\right)^{1 / p^{\prime}}=B, \quad\left(\int f^{p} d x\right)^{1 / p}=F, \quad\left(\int g^{p^{\prime}} d y\right)^{1 / p^{\prime}}=G$, and (iv) $p>1$, then

$$
\begin{gather*}
\Sigma \Sigma \frac{a_{m} b_{n}}{\operatorname{Max}(m, n)}<p p^{\prime} A B  \tag{1}\\
\iint \frac{f(x) g(y)}{\operatorname{Max}(x, y)} d x d y<p p^{\prime} F G \tag{2}
\end{gather*}
$$

unless $\left(a_{m}\right)$ or $\left(b_{n}\right)$ or $f(x)$ or $g(y)$ is nul. The constants are the best possible.
[Cases of Theorems 318 and 319, (a). In order to shorten the statements of the following theorems we agree that conditions (i), (ii), and (iii) are presupposed in all of them; and that, whenever the conclusion is expressed by an inequality

$$
X<K Y \quad(\text { or } X>K Y)
$$

with a definite $K$, then $K$ has its best possible value (unless the contrary is stated explicitly) and equality is excluded unless a sequence or function involved in the theorem is nul.

When, on the other hand, the conclusion is

$$
X \leqq K Y
$$

with an unspecified $K$, then $K$ is a function of any parameters of the theorem.]
342. If $p>1$, then

$$
\begin{gather*}
\Sigma \Sigma \frac{\log (m / n)}{m-n} a_{m} b_{n}<\pi^{2} \operatorname{cosec}^{2} \frac{\pi}{p} \cdot A B  \tag{1}\\
\iint \frac{\log (x / y)}{x-y} f(x) g(y) d x d y<\pi^{2} \operatorname{cosec}^{2} \frac{\pi}{p} \cdot F G
\end{gather*}
$$

[Also cases of Theorems 318 and 319, (a). Here

$$
\left.k=\int_{0}^{\infty} \frac{\log x}{x-1} x^{-1 / p} d x=\pi^{2} \operatorname{cosec}^{2} \frac{\pi}{p} .\right]
$$

343. If $p>1$, then

$$
\sum_{2}^{\infty} \sum_{2}^{\infty} \frac{a_{m} b_{n}}{m n \log m n}<\frac{\pi}{\sin (\pi / p)}\left(\sum_{2}^{\infty} \frac{a_{m}{ }^{p}}{m}\right)^{1 / p} \sum_{2}^{\infty}\left(\frac{b_{n}{ }^{p \prime}}{n}\right)^{1 / p^{\prime}} .
$$

[Mulholland (2). Since

$$
\log \frac{m+1}{m}<\frac{1}{m}
$$

the result is slightly stronger than that obtained by taking $\Lambda_{m}=\log m$, $\mathrm{M}_{n}=\log n$ in Theorem 321.]
344. If $0<p<1$ then

$$
\Sigma\left(a_{n}+a_{n+1}+\ldots\right)^{p}>p^{p} \Sigma\left(n a_{n}\right)^{p} .
$$

[Copson (2). This theorem, with Theorems 326, 331, and 338, forms a systematic set of four.]
345. If $0<p<1$ then

$$
\Sigma\left(\frac{a_{n}+a_{n+1}+\ldots}{n}\right)^{p}>p^{p} \Sigma a_{n}{ }^{p} .
$$

[Corollary of Theorem 344. Compare Theorem 338; here there is no gloss, but the constant is less favourable and is presumably not the best possible.]
346. If (a) $c>1, s_{n}=a_{1}+a_{2}+\ldots+a_{n}$, or (b) $c<1, s_{n}=a_{n}+a_{n+1}+\ldots$, then
( $\alpha$ )

$$
\begin{array}{ll}
\Sigma n^{-c} s_{n}^{p} \leqq K \Sigma n^{-c}\left(n a_{n}\right)^{p} & (p>1), \\
\Sigma n^{-c} s_{n}^{p} \geqq K \Sigma n^{-c}\left(n a_{n}\right)^{p} & (0<p<1) .
\end{array}
$$

[In each of the four cases $K=K(p, c)$, as laid down under Theorem 341. See Hardy and Littlewood (1).

We prove $(\alpha)$ when $c>1$. If

$$
\phi_{n}=n^{-c}+(n+1)^{-c}+\ldots,
$$

then $\phi_{n}<K n^{1-c}$. Hence, if we agree that $s_{0}=0$, we have

$$
\begin{aligned}
\sum_{1}^{m} n^{-c} s_{n}^{p} & =\sum_{1}^{m}\left(\phi_{n}-\phi_{n+1}\right) s_{n}^{p} \leqq \sum_{1}^{m} \phi_{n}\left(s_{n}^{p}-s_{n-1}{ }^{p}\right) \\
& \leqq K \sum_{1}^{m} n^{1-c} s_{n}^{p-1} a_{n} \leqq K\left(\sum_{1}^{m} n^{-c}\left(n a_{n}\right)^{p}\right)^{1 / p}\left(\sum_{1}^{m} n^{-c} s_{n}^{p}\right)^{1 / p^{\prime}} ;
\end{aligned}
$$

and ( $\alpha$ ) follows. Hardy and Littlewood (2) give function-theoretic applications of $(\alpha)$ and $(\beta)$. The important case is that in which $c=2$.]
347. If $r$ and $F$ satisfy the conditions of Theorem 330, but $0<p<1$, then

$$
\int x^{-r} F^{v} d x>\left(\frac{p}{|r-1|}\right)^{p} \int x^{-r}(x f)^{p} d x
$$

[Hardy (5).]
348. If

$$
\sigma(y)=\Sigma a_{m} e^{-m / y}
$$

and $p>1$, then

$$
2^{-p} \Sigma\left\{\frac{\sigma(n)}{n}\right\}^{p}<\int_{0}^{\infty}\left\{\frac{\sigma(y)}{y}\right\}^{p} d y<\left\{\Gamma\left(\frac{1}{p^{\prime}}\right)\right\}^{p} A^{p}
$$

[Take $K(x, y)=y^{-1} e^{-x / y}$ and apply Theorem 319, (b). More general but less precise results are given by Hardy and Littlewood (1), and some function-theoretic applications by Hardy and Littlewood (2).]
349. If $\lambda_{n}$ and $\Lambda_{n}$ satisfy the conditions of Theorem 332, then

$$
\Sigma \lambda_{n}\left(a_{1} \lambda_{1} a_{2}^{\lambda_{2}} \ldots a_{n}^{\lambda_{n}}\right)^{1 / \Lambda_{n}}<e \Sigma \lambda_{n} a_{n} .
$$

[See Hardy (4).]
350. If $p>1, K(x)>0$, and

$$
\int K(x) x^{s-1} d x=\phi(s),
$$

then

$$
\begin{gathered}
\iint K(x y) f(x) g(y) d x d y<\phi\left(\frac{1}{p}\right)\left(\int x^{p-2} f^{p} d x\right)^{1 / p}\left(\int g^{p^{\prime}} d y\right)^{1 / p^{\prime}} \\
\int d x\left(\int K(x y) f(y) d y\right)^{p}<\phi^{p}\left(\frac{1}{p}\right) \int x^{p-2} f^{p} d x \\
\int x^{p-2} d x\left(\int K(x y) f(y) d y\right)^{p}<\phi^{p}\left(\frac{1}{p^{\prime}}\right) \int f^{p} d x
\end{gathered}
$$

In particular, when $K(x)=e^{-x}$, and $F(x)=\int K(x y) f(y) d y$ is the 'Laplace transform' of $f(x)$,

$$
\int F^{y} d x<\Gamma^{p}\left(\frac{1}{p}\right) \int x^{p-2} f^{p} d x, \quad \int x^{p-2} F^{p} d x<\Gamma^{p}\left(\frac{1}{p^{\prime}}\right) \int f^{p} d x
$$

351. If also $K(x)$ is a decreasing function of $x$, and

$$
A(x)=\Sigma a_{n} K(n x), \quad A_{n}=\int a(x) K(n x) d x,
$$

then

$$
\begin{gathered}
\int A^{p}(x) d x<\phi^{p}\left(\frac{\mathbf{l}}{p}\right) \Sigma n^{p-2} a_{n}^{p}, \\
\Sigma A_{n}^{p}<\phi^{p}\left(\frac{\mathbf{l}}{p}\right) \int x^{p-2} a^{p}(x) d x, \\
\int x^{p-2} A^{p}(x) d x<\phi^{p}\left(\frac{\mathbf{1}}{p^{\prime}}\right) \Sigma a_{n}^{p}, \\
\Sigma n^{p-2} A_{n}^{p}<\phi^{p}\left(\frac{\mathbf{1}}{p^{\prime}}\right) \int a^{p}(x) d x .
\end{gathered}
$$

352. If $F^{\prime}(x)$ is the Laplace transform of $f(x)$, and $1<p \leqq 2$, then

$$
\int F^{p^{\prime}} d x \leqq \frac{2 \pi}{p^{\prime}}\left(\int f^{p} d x\right)^{p^{\prime} / p}
$$

[For the last three theorems see Hardy (10). Theorem 350 may be deduced from Theorem 319 by transformation. It is not asserted that the constant in Theorem 352 is the best possible.]
353. If $K_{0}(x) \geqq 0$,

$$
\begin{gathered}
K_{1}(x, y)=\int K_{0}(x t) K_{0}(y t) d t, \quad K_{2}(x, y)=\int K_{1}(x, t) K_{1}(y, t) d t \\
\int \frac{K_{1}(x, 1)}{\sqrt{ } x} d x=k
\end{gathered}
$$

then

$$
\Sigma \Sigma K_{2}(m, n) a_{m} a_{n} \leqq k \Sigma \Sigma K_{1}(m, n) a_{m} a_{n}
$$

[See Hardy (9). The theorem is one concerning quadratic, not bilinear, forms.]
354.

$$
\Sigma \Sigma \frac{\log (m / n)}{m-n} a_{m} a_{n} \leqq \pi \Sigma \Sigma \frac{a_{m} a_{n}}{m+n}
$$

355. 

$$
\Sigma \Sigma \frac{|\log (m / n)|}{\operatorname{Max}(m, n)} a_{m} a_{n} \leqq 2 \Sigma \Sigma \frac{a_{m} a_{n}}{\operatorname{Max}(m, n)}
$$

[Corollaries of Theorem 353. Observe that Theorem 354, when combined with Theorem 315, gives

$$
\Sigma \Sigma \frac{\log (m / n)}{m-n} a_{m} a_{n} \leqq \pi^{2} \Sigma a_{m}^{2}
$$

in agreement with Theorem 342.]
356. If

$$
c(x)=\int_{0}^{x} a(t) b(x-t) d t,
$$

$A^{p}=\int x^{-1}\left(x^{a} a(x)\right)_{a}^{p} d x, \quad B^{q}=\int x^{-1}\left(x^{\beta} b(x)\right)^{q} d x, \quad C^{r}=\int x^{-1}\left(x^{\gamma} c(x)\right)^{r} d x$.

$$
p>1, \quad q>1, \quad \frac{1}{r} \leqq \frac{1}{p}+\frac{1}{q}, \quad \alpha<1, \quad \beta<1, \quad \gamma=\alpha+\beta-1,
$$

then

$$
C<K A B
$$

$$
K=\frac{\Gamma(1-\alpha) \Gamma(1-\beta)}{\Gamma(1-\gamma)} .
$$

357. If $a_{0}=b_{0}=0, c_{n}=a_{0} b_{n}+a_{1} b_{n-1}+\ldots+a_{n} b_{0}$,
$A^{p}=\Sigma n^{-1}\left(n^{a} a_{n}\right)^{p}, \quad B^{q}=\Sigma n^{-1}\left(n^{\beta} b_{n}\right)^{q}, \quad C^{r}=\Sigma n^{-1}\left(n^{\gamma} c_{n}\right)^{r}$,
$p, q, r$, and $\gamma$ satisfy the conditions of Theorem 356, and $0 \leqq \alpha<1$, $0 \leqq \beta<1$, then $C<K A B$, with the $K$ of Theorem 356. If $\alpha<0$ or $\beta<0$, the inequality is true with some $K$.
358. If

$$
a_{0}=b_{0}=\ldots=c_{0}=0
$$

and

$$
u_{n}=\Sigma a_{r_{1}} b_{r_{2}} \ldots c_{r_{k}} \quad\left(r_{j} \geqq 0, \Sigma r_{j}=n\right)
$$

then $\quad \Sigma u_{n}{ }^{2}<\frac{1}{\pi}\left\{\Gamma\left(\frac{1}{2 k}\right)\right\}^{2 k}\left(\Sigma n^{2 k-2} a_{n}^{2 k}\right)^{1 / k} \ldots\left(\Sigma n^{2 k-2} c_{n}^{2 k}\right)^{1 / k}$.
359. If $p>1, l>0, m>0$ and $c(x)$ is defined as in Theorem 356, then

$$
\int x^{(1-l-m)(p-1)} c^{p}(x) d x \leqq K \int x^{(1-l)(p-1)} a^{p}(x) d x \int x^{(1-m)(p-1)} b^{p}(x) d x
$$

where

$$
K=\left\{\frac{\Gamma(l) \Gamma(m)}{\Gamma(l+m)}\right\}^{p-1}
$$

There is equality if and only if

$$
a(x) \equiv A x^{l-1} e^{-C x}, \quad b(x) \equiv B x^{m-1} e^{-C x},
$$

where $A, B, C$ are non-negative constants and $C$ is positive.
[For Theorems 356-359 see Hardy and Littlewood (3, 5, and 12).]
360. If $L(x)$ is the Laplace transform of $f(x)$, and $q \geqq p>1$, then

$$
\int x^{-(p+q-p q) / p} L^{q}(x) d x \leqq K F^{q}
$$

361. If $p>1, q>1$,

$$
\mu=\frac{1}{p}+\frac{1}{q}-1 \geqq 0
$$

and $L, M$ are the Laplace transforms of $f, g$, then

$$
\int x^{-\mu} L M d x \leqq K F G
$$

362. If $p>1,0 \leqq \mu<1 / p$, and
then

$$
\begin{gathered}
\alpha_{n}=\sum_{m} \frac{a_{m}}{(m+n)^{1-\mu}}, \\
\Sigma \alpha_{n}{ }^{p /(1-\mu p)} \leqq K A^{p /(1-\mu p)} .
\end{gathered}
$$

[This may be deduced from Theorem 339 by the converse of Hölder's inequality. Many further theorems of the same general character as Theorems 360-362 are given by Hardy and Littlewood (1).]
363. If $\lambda_{n}$ is positive and

$$
p>1, \quad A_{n}=a_{1}+a_{2}+\ldots+a_{n}, \quad \lambda_{1}+\lambda_{2}+\ldots+\lambda_{n} \leqq c n,
$$

then

$$
\Sigma\left(\frac{A_{n}}{n}\right)^{p} \lambda_{n} \leqq K A^{p}
$$

364. If $\lambda_{n}$ is positive and

$$
\begin{gathered}
p>1, \quad r>1, \quad \lambda_{1}^{r}+\lambda_{2}^{r}+\ldots+\lambda_{n}^{r} \leqq c n, \\
\Sigma \Sigma \frac{a_{m} b_{n}}{m+n} \lambda_{m+n} \leqq K A B .
\end{gathered}
$$

then
The result is not necessarily true when $r=1$ (as it is when $\lambda_{n}=1$ ).
[For these two theorems, which are corresponding extensions of Theorems 326 and 315, see Hardy and Littlewood (11).]
365. The inequalities in Theorems 326 and 334 are the special cases $\phi=x^{t}(0<t<1), \phi=\log x$, of

$$
\begin{equation*}
\Sigma \phi^{-1}\left(\frac{\phi\left(a_{1}\right)+\ldots+\phi\left(a_{n}\right)}{n}\right)<K(\phi) \Sigma a_{n} . \tag{i}
\end{equation*}
$$

[Knopp (2). This remark has led Knopp to a systematic investigation of forms of $\phi$ for which (i) is true. See also Mulholland (4).]
366. Suppose that $\phi$ and $\psi$ are continuous and strictly increasing for $x>0$, and have the limits 0 or $-\infty$ when $x \rightarrow 0$; and that $\phi$ is convex with respect to $\psi(\S 3.9)$. Then (i), if true for $\phi$, is also true for $\psi$, with $K(\psi) \leqq K(\phi)$.
[Knopp (2).]
367. $\quad \Sigma\left(\log \frac{e^{1 / a_{1}}+e^{1 / a_{2}}+\ldots+e^{1 / a_{n}}}{n}\right)^{-1}<2 \Sigma a_{n}$.
[Knopp (2).]

## CHAPTER X

## REARRANGEMENTS

10.1. Rearrangements of finite sets of variables. In what follows we are concerned with finite sets of non-negative numbers such as

$$
\begin{gathered}
a_{1}, a_{2}, \ldots, a_{j}, \ldots, a_{n} ; \quad b_{1}, b_{2}, \ldots, b_{j}, \ldots, b_{n} \\
a_{-n}, \ldots, a_{0}, \ldots, a_{j}, \ldots, a_{n}
\end{gathered}
$$

we denote such sets by $(a),(b), \ldots$.
Taking for example the first set, in which $j$ assumes the values $1,2, \ldots, n$, we define a permutation function $\phi(j)$ as a function which takes each of the values $1,2, \ldots, n$ just once when $j$ varies through the same aggregate of values. If

$$
a_{\phi(j)}=a_{j}^{\prime} \quad(j=1,2, \ldots, n)
$$

then we describe ( $a^{\prime}$ ) as a rearrangement of ( $a$ ). Similar definitions apply to other cases in which the range of variation of $j$ is different.

There are certain special rearrangements of (a) which are particularly important here. These rearrangements, which we denote by

$$
(\bar{a}),\left(a^{+}\right),(+a),\left(a^{*}\right)
$$

are defined as follows.
The set $(\bar{a})$ is the set (a) rearranged in ascending order, so that, when the values of $j$ are $1,2, \ldots, n$,

$$
\bar{a}_{1} \leqq \bar{a}_{2} \leqq \ldots \leqq \bar{a}_{n} .
$$

The set ( $\bar{a}$ ) is defined unambiguously by the set ( $a$ ) although, when the $a$ are not all different, there are ambiguities in the definition of the permutation function by which we pass from (a) to ( $\bar{a}$ ).

In defining the sets $\left(a^{+}\right),\left({ }^{+} a\right),\left(a^{*}\right)$ we suppose that $j$ varies from $-n$ to $n$. The set ( $a^{+}$) is defined by

$$
a_{0}{ }^{+} \geqq a_{1}{ }^{+} \geqq a_{-1}{ }^{+} \geqq a_{2}{ }^{+} \geqq a_{-2}{ }^{+} \geqq \ldots
$$

and the set ( ${ }^{+} a$ ) by

$$
+a_{0} \geqq+a_{-1} \geqq+a_{1} \geqq+a_{-2} \geqq+a_{2} \geqq \ldots .
$$

There is one particularly important case, that in which every value of an $a$, except the largest, occurs an even number of times, while the largest value occurs an odd number of times. In this case we shall say that the set $(a)$ is symmetrical. The sets $\left(a^{+}\right)$and $(+a)$ are then identical, and we write

$$
a^{+}=+a=a^{*}
$$

so that $a^{*}$ is defined by

$$
a_{0}^{*} \geqq a_{1}^{*}=a_{-1}^{*} \geqq a_{2}^{*}=a_{-2}^{*} \geqq \ldots
$$

A set ( $a^{*}$ ) may be said to be symmetrically decreasing. The sets $\left(a^{+}\right)$and $\left({ }^{+} a\right)$ are sets arranged so as to be as nearly symmetrically decreasing as possible, but with the inevitable overweight of one side arranged systematically to the advantage of the right or the left respectively. All these sets are defined unambiguously by ( $a$ ), though there may be ambiguities in the definitions of the corresponding permutation functions.

We note that

$$
\begin{equation*}
a_{j}^{+}=+a_{-j} \tag{10.1.1}
\end{equation*}
$$

10.2. A theorem concerning the rearrangements of two sets. We begin by proving a verysimple, but important, theorem concerning the set $(\bar{a})$.
368. ${ }^{\text {a }}$ If (a) and (b) are given except in arrangement, then

$$
\Sigma a b
$$

is greatest when (a) and (b) are monotonic in the same sense and least when they are monotonic in opposite senses; that is to say

$$
\begin{equation*}
\sum_{j=1}^{n} \bar{a}_{j} \bar{b}_{n+1-j} \leqq \sum_{1}^{n} a_{j} b_{j} \leqq \sum_{1}^{n} \bar{a}_{j} \bar{b}_{j} \tag{10.2.1}
\end{equation*}
$$

It will be observed that, since we can add up the sum $\Sigma a b$ in any order, we may suppose one set, say (a), arranged from the beginning in any order we please (in particular in ascending order).

We may express the theorem equally well by saying that the maximum corresponds to 'similar ordering' of $(a)$ and $(b)$ in the

[^68]sense of $\S 2.17$, the minimum to 'opposite ordering'a. The theorem becomes 'intuitive' if we interpret the $a$ as distances along a rod to hooks and the $b$ as weights suspended from the hooks. To get the maximum statical moment with respect to an end of the rod, we hang the heaviest weights on the hooks farthest from that end.

To prove the theorem, suppose that the ( $a$ ) are in ascending order, but not the (b). Then there are a $j$ and a $k$ such that $a_{j} \leqq a_{k}$ and $b_{j}>b_{k}$. Since

$$
a_{j} b_{k}+a_{k} b_{j}-\left(a_{j} b_{j}+a_{k} b_{k}\right)=\left(a_{k}-a_{j}\right)\left(b_{j}-b_{k}\right) \geqq 0
$$

we do not diminish $\Sigma a b$ by exchanging $b_{j}$ and $b_{k}$. A finite number of such exchanges leads to an ascending order of the $b$, so that

$$
\Sigma a b \leqq \Sigma \bar{a} \bar{b}
$$

The other half of the theorem is proved in the same way.
This argument establishes incidentally a variant of Theorem 368 which is sometimes useful.
369. If
(10.2.2)

$$
\Sigma a b^{\prime} \leqq \Sigma a b
$$

for all rearrangements $\left(b^{\prime}\right)$ of $(b)$, then $(a)$ and $(b)$ are similarly ordered.

For, if $\left(a_{j}-a_{k}\right)\left(b_{j}-b_{k}\right)<0$ for any $j, k$, we can falsify (10.2.2) by exchanging $b_{j}$ and $b_{k}$.
10.3. A second proof of Theorem 368. We have to consider analogues of Theorem 368 for more than two sets of variables. These lie a good deal deeper and cannot be proved in so simple a manner. We therefore give a second proof of Theorem 368 which, though quite unnecessarily complicated for its immediate object, will serve to introduce the method which we use later. We confine ourselves to the second inequality (10.2.1), and divide the proof into three stages.
(1) Suppose first that the sets considered consist entirely of

[^69]0 's and 1's; we indicate such special sets by the use of German letters $\mathfrak{a}, \mathfrak{b}, \ldots$ Then

$$
\begin{equation*}
\mathfrak{a}^{2}=\mathfrak{a}, \quad \mathfrak{b}^{2}=\mathfrak{b} \tag{10.3.1}
\end{equation*}
$$

for all $j$. In this case

$$
\Sigma \mathfrak{a b} \leqq \Sigma \mathfrak{a}, \quad \Sigma \mathfrak{a b} \leqq \Sigma \mathfrak{b}
$$

and so $\quad \Sigma \mathfrak{a} \mathfrak{b} \leqq \operatorname{Min}(\Sigma \mathfrak{a}, \Sigma \mathfrak{b})=\Sigma \overline{\mathfrak{a}} \overline{\mathfrak{b}}$.
(2) Any set (a) may be decomposed into a linear combination of sets

$$
\left(\mathfrak{a}^{1}\right), \quad\left(\mathfrak{a}^{2}\right), \quad \ldots, \quad\left(\mathfrak{a}^{l}\right)^{\mathfrak{a}}
$$

of the special type considered under (1), in such a way that

$$
\begin{equation*}
a_{j}=\alpha^{1} a_{j}^{1}+\alpha^{2} a_{j}^{2}+\ldots+\alpha^{l} \mathfrak{a}_{j}^{l} \quad(j=1,2, \ldots, n) \tag{10.3.2}
\end{equation*}
$$

and
(10.3.3) $\quad \bar{a}_{j}=\alpha^{1} \overline{\mathfrak{a}}_{j}^{1}+\alpha^{2} \overline{\mathfrak{a}}_{j}^{2}+\ldots+\alpha^{l} \overline{\mathfrak{a}}_{j}^{l} \quad(j=1,2, \ldots, n)$,
the coefficients $\alpha$ being non-negative.
The method of decomposition will become clear by considering a special case. Suppose that ( $a$ ) contains (in some order) the three numbers $A, B, C$, where $0 \leqq A \leqq B \leqq C$, so that

$$
\bar{a}_{1}=A, \quad \bar{a}_{2}=B, \quad \bar{a}_{3}=C .
$$

Then

$$
\begin{aligned}
& \bar{a}_{1}=A .1+(B-A) 0+(C-B) 0, \\
& \bar{a}_{2}=A .1+(B-A) 1+(C-B) 0, \\
& \bar{a}_{3}=A .1+(B-A) 1+(C-B) 1
\end{aligned}
$$

and we may write

$$
\begin{gathered}
\bar{a}_{j}=\alpha^{1} \overline{\mathfrak{a}}_{j}^{1}+\alpha^{2} \overline{\mathfrak{a}}_{j}^{2}+\alpha^{3} \overline{\mathfrak{a}}_{j}^{3}, \\
\alpha^{1}=A, \quad \alpha^{2}=B-A, \quad \alpha^{3}=C-B
\end{gathered}
$$

where
and $\left(\overline{\mathfrak{a}}^{1}\right),\left(\overline{\mathfrak{a}}^{2}\right),\left(\overline{\mathfrak{a}}^{3}\right)$ are the three sets

$$
(1,1,1), \quad(0,1,1), \quad(0,0,1)
$$

If then we perform the permutation which changes $(\bar{a})$ into ( $a$ ) and at the same time $\left(\overline{\mathfrak{a}}^{1}\right), \ldots$ into $\left(\mathfrak{a}^{1}\right), \ldots,{ }^{b}$ we obtain

$$
a_{j}=\alpha^{1} \mathfrak{a}_{j}^{1}+\alpha^{2} \mathfrak{a}_{j}^{2}+\alpha^{3} \mathfrak{a}_{j}^{3}
$$

[^70]In the general case we proceed in the same way, writing

$$
\begin{aligned}
& \bar{a}_{1}=\bar{a}_{1} \cdot 1+\left(\bar{a}_{2}-\bar{a}_{1}\right) \cdot 0+\left(\bar{a}_{3}-\bar{a}_{2}\right) \cdot 0+\ldots, \\
& \bar{a}_{2}=\bar{a}_{1} \cdot 1+\left(\bar{a}_{2}-\bar{a}_{1}\right) \cdot 1+\left(\bar{a}_{3}-\bar{a}_{2}\right) \cdot 0+\ldots,
\end{aligned}
$$

This secures (10.3.3), and (10.3.2) then follows by rearrangement, as in the special case.
(3) From (1) and (2) we can deduce the general theorem. For, decomposing (b) as in (2), we have

$$
\begin{aligned}
& a_{j}=\sum_{\rho} \alpha^{\rho} \mathfrak{a}_{j}^{\rho}, \quad \bar{a}_{j}=\sum_{\rho} \alpha^{\rho} \overline{\mathfrak{a}}_{j}^{\rho}, \quad b_{j}=\sum_{\sigma} \beta^{\sigma} \mathfrak{b}_{j}^{\sigma}, \quad \bar{b}_{j}=\sum_{\sigma} \beta^{\sigma} \overline{\mathfrak{b}}_{j}^{\sigma}, \\
& \sum_{j} a_{j} b_{j}=\sum_{\rho \sigma} \sum^{\rho} \alpha^{\rho} \beta^{\sigma} \sum_{j} \mathfrak{a}_{j}^{\rho} \mathfrak{b}_{j}^{\sigma} \\
& \leqq \sum_{\rho \sigma} \alpha^{\rho} \beta^{\sigma} \sum_{j} \overline{\mathfrak{a}}_{j} \rho \overline{\mathfrak{b}}_{j}^{\sigma}=\sum_{j} \bar{a}_{j} \bar{b}_{j} .
\end{aligned}
$$

10.4. Restatement of Theorem 368. It will also be useful to restate Theorem 368 in different language. We suppose now that, in the sets $(a),(b), j$ runs from $-n$ to $n$. We write

$$
f(x)=\Sigma a_{j} x^{j}, \quad g(x)=\Sigma b_{j} x^{j}
$$

and call

$$
a_{0}=\mathfrak{C}(f(x))
$$

the central coefficient of $f$. Plainly

$$
\mathfrak{C}\left(f\left(x^{-1}\right)\right)=\mathfrak{C}(f(x)) .
$$

Also

$$
\sum_{r+s=0} a_{r} b_{s}=\Sigma a_{j} b_{-j}=\mathfrak{C}(f g)
$$

The sets $\left(a_{j}{ }^{+}\right)$and $\left({ }^{+} b_{-j}\right)$ are similarly ordered, and if we write

$$
f^{+}(x)=\Sigma a_{j}+x^{j}, \quad+f(x)=\Sigma+a_{j} x^{j},
$$

so that, by (10.1.1),

$$
f^{+}\left(x^{-1}\right)=+f(x)
$$

then Theorem 368 gives

$$
\mathfrak{C}(f g)=\sum_{r+s=0} a_{r} b_{s}=\Sigma a_{j} b_{-j} \leqq \Sigma a_{j}^{++} b_{-j}=\sum_{r+s=0} a_{r}^{++} b_{s}=\mathfrak{C}\left(f^{++} g\right) .
$$

Hence we deduce
370. The central coefficient of

$$
\sum_{-n}^{n} a_{i} x^{i} \sum_{-n}^{n} b_{j} x^{j}
$$

is greatest, for all rearrangements of the $a$ and $b$, when $\left(a_{j}\right)$ and $\left(b_{-j}\right)$ are similarly ordered, in particular when $(a)$ is $\left(a^{+}\right)$and $(b)$ is $(+b)$ or when $(a)$ is $\left({ }^{+} a\right)$ and $(b)$ is $\left(b^{+}\right)$.
10.5. Theorems concerning the rearrangements of three sets. We pass now to theorems involving three sets of variables.
371.a Suppose that the $c, x$, and $y$ are non-negative, and the $c$ symmetrically decreasing, so that

$$
c_{0} \geqq c_{1}=c_{-1} \geqq c_{2}=c_{-2} \geqq \ldots \geqq c_{2 k}=c_{-2 k},
$$

while the $x$ and $y$ are given except in arrangement. Then the bilinear form

$$
S=\sum_{r=-k}^{k} \sum_{s=-k}^{k} c_{r-s} x_{r} y_{s}
$$

attains its maximum when $(x)$ is $\left(x^{+}\right)$and $y$ is $\left(y^{+}\right)$.
It is evident that, if this is so, then the maximum must also be attained when $(x)$ is $\left({ }^{+} x\right)$ and $(y)$ is $(+y)$.
372.b Suppose that (a), (b), (c) are three sets satisfying
(10.5.1) $\quad a_{0} \geqq a_{r}=a_{-r}, \quad b_{0} \geqq b_{s}=b_{-s}, \quad c_{0} \geqq c_{t}=c_{-i}$.

Then the maximum of

$$
\sum_{r+s+t=0} a_{r} b_{s} c_{t}=\mathfrak{C}\left(\Sigma a_{r} x^{r} \Sigma b_{s} x^{s} \Sigma c_{t} x^{t}\right),
$$

for rearrangements of the sets which leave $a_{0}, b_{0}, c_{0}$ unaltered, is attained when (a), (b), (c) are ( $a^{*}$ ), ( $b^{*}$ ), ( $\left.c^{*}\right)$.
373.c If (a), (b), (c) are three sets, of which (c) is symmetrical in the sense of $\S 10.1$, then

$$
\sum_{r+s+t=0} a_{r} b_{s} c_{t} \leqq \sum_{r+s+t=0} a_{r}++b_{s} c_{t}^{*}=\sum_{r+s+t=0}+a_{r} b_{s}+c_{t}{ }^{*} .
$$

It will be sufficient to prove Theorem 373, since this includes theother two theorems. In the first place, Theorem 373 is Theorem 372 freed from the restrictions (10.5.1), wholly in regard to ( $a$ ) and ( $b$ ) and partly in regard to (c). To deduce Theorem 371 from Theorem 373, we put $2 k=n, x_{r}=a_{-r}, y_{s}=b_{s}$, and suppose that the $a$ and $b$ outside the range $(-k, k)$ are zero. We may observe

[^71]finally that Theorem 370 is the simple case of Theorem 373 in which $c_{0}=1$ and the remaining $c$ are 0 .
10.6. Reduction of Theorem 373 to a special case. We divide the proof of Theorem 373 into three stages, as in $\S 10.3$. The whole difficulty of the proof lies in stage (1), in which (a), (b), (c) are of types $(\mathfrak{a}),(\mathfrak{b}),(\mathfrak{c})$ : and we take this stage for granted for the moment and dispose of the easier stages (2) and (3).

First, we may decompose (a), (b), (c) into sums of sets ( $\mathfrak{a}^{\rho}$ ), $\left(\mathfrak{b}^{\sigma}\right),\left(\mathfrak{c}^{\tau}\right)$, in such a way that

$$
a_{j}=\sum_{\rho} \alpha^{\rho} \mathfrak{a}_{j}{ }^{\rho}, \quad b_{j}=\underset{\sigma}{\Sigma} \beta^{\sigma} \mathfrak{b}_{j}{ }^{\sigma}, \quad c_{j}=\sum_{\tau} \gamma^{\tau} c_{j}^{\tau},
$$

and

$$
a_{j}{ }^{+}=\sum_{\rho} \alpha^{\rho} \mathfrak{a}_{j}{ }^{\rho+}, \quad+b_{j}=\underset{\sigma}{\Sigma} \beta^{\sigma}+\mathfrak{b}_{j}{ }^{\sigma}, \quad c_{j}{ }^{*}=\underset{\tau}{\Sigma_{\gamma} \mathcal{c}_{j}{ }^{\tau *} .}
$$

Here the $\mathfrak{a}, \mathfrak{b}, \mathfrak{c}$ are all 0 or l , the $\alpha, \beta, \gamma$ are non-negative, and (a point which doesnotarise in $\S 10.3$ ) the sets $\left(\mathfrak{c}^{\tau}\right)$ are symmetrical. Allthis is proved by the method of $\S 10.3(2) .^{a}$ When we have done this, and proved the theorem for sets of type (a), (b), ( $\mathfrak{c}$ ), we have

$$
\begin{aligned}
\sum_{r+s+t=0}^{\sum} a_{r} b_{s} c_{t}= & =\sum_{\rho, \sigma, \tau} \alpha^{\rho} \beta^{\sigma} \gamma^{\tau}{ }_{r+s+t=0}^{\Sigma} \mathfrak{a}_{r}{ }^{\rho} \mathfrak{b}_{s}{ }^{\sigma} \mathfrak{c}_{t}^{\tau} \\
& \leqq \sum_{\rho, \sigma, \tau} \alpha^{\rho} \beta^{\sigma} \gamma^{\tau} \sum_{r+s+t=0} \mathfrak{a}_{r}^{\rho++} \mathfrak{b}_{s}{ }^{\sigma} \mathfrak{c}_{t}^{\tau *}=\underset{r+s+t=0}{\sum} a_{r}++b_{s} c_{t}^{*},
\end{aligned}
$$

and the proof is completed.
It remains to prove the theorem in the special case in which all $a, b, c$ are 0 or 1.b The set $c$, being symmetrical, contains an even number of 0 's and an odd number of 1 's. We write

$$
f(x)=\Sigma a_{r} x^{r}, \quad g(x)=\Sigma b_{s} x^{s}, \quad h(x)=\Sigma c_{t} x^{x^{t}} .
$$

Since we may add any number of 0 's to the sets, we may suppose that all the summations run from $-n$ to $n$.

We have also

$$
\begin{aligned}
& f^{+}(x)=\Sigma a_{r}+x^{r}=x^{-R}+\ldots+1+\ldots+x^{R^{\prime}} \\
& +g(x)=\Sigma+b_{s} x^{s}=x^{-S^{\prime}}+\ldots+1+\ldots+x^{S} \\
& h^{*}(x)=\Sigma c_{l}^{*} x^{t}=x^{-T}+\ldots+1+\ldots+x^{T}
\end{aligned}
$$

[^72]where $R, R^{\prime}, S, S^{\prime}, T$ are non-negative integers and
(10.6.1) $\quad R \leqq R^{\prime} \leqq R+1, \quad S \leqq S^{\prime} \leqq S+1$.

We have to prove that

$$
\begin{equation*}
\mathfrak{C}(f g h) \leqq \subseteq\left(f^{++} g h^{*}\right) . \tag{10.6.2}
\end{equation*}
$$

The inequality (10.6.2) may be made 'intuitive' by a geometrical representation. Let $x, y$ be rectangular coordinates in a plane, and represent each non-zero coefficient of $f, g, h$ by a line, $x=r$ for $a_{r}=1, y=s$ for $b_{s}=1$, and $x+y=-t$ for $c_{t}=1$. If $a_{r} b_{s} c_{t}$ contributes a unit to $\mathfrak{C}(f g h)$, these three lines intersect. Each of the functions $f, g, h$ is represented by a family of parallel lines,


Fig. 1. Graph of $f, g, h$.


Fig. 2. Graph of $f^{+},+g, h^{*}$.
and $\mathfrak{C}(f g h)$ is the total number of triple intersections of these lines. We represent $f^{+},+g, h^{*}$ similarly; $f^{+}$is also represented by $R+1+R^{\prime}$ vertical lines, but now these lines are shifted as near as possible together. Typical figures are shown in Figs. 1 and 2: here $(a),(b),(c)$ are the sets

$$
\begin{gathered}
1,0,1,0,0,1,0,0,1 \\
1,1,0,0,1,1,1,0,0,1 \\
1,0,1,0,0,0,1,0,1,0,0,1,0
\end{gathered}
$$

respectively; and

$$
R=1, \quad R^{\prime}=2, \quad S=2, \quad S^{\prime}=3, \quad T=2
$$

It is intuitive that the number of intersections is greatest when, as in Fig. 2, the diagram is as condensed as possible.

Our proof of (10.6.2) may be presented geometrically ${ }^{a}$ and followed on the figures. We reduce the actual case considered to a simpler one by taking away one horizontal and one vertical line from each figure, as is suggested by the thick lines in the figures. We prefer however to state the proof in a purely analytical form.
10.7. Completion of the proof. There are three subordinate cases in which the proof is easy.
(1) If $R^{\prime}=0, f^{+}$reduces to 1 and the result is included in Theorem 370.
(2) If $S^{\prime}=0,+g$ reduces to 1 and again the result is included in Theorem 370.
(3) Suppose that

$$
\begin{equation*}
R+S^{\prime} \leqq T, \quad R^{\prime}+S \leqq T \tag{10.7.1}
\end{equation*}
$$

We have in any case

$$
\begin{equation*}
\mathfrak{C}(f g h)=\sum_{r+s+t=0} a_{r} b_{s} c_{t} \leqq \Sigma a_{r} \Sigma b_{s}=\left(R+1+R^{\prime}\right)\left(S+1+S^{\prime}\right) \tag{10.7.2}
\end{equation*}
$$

But, when the inequalities (10.7.1) are satisfied,
$\mathfrak{C}\left(f^{++} g h^{*}\right)=\mathfrak{C}\left\{\left(x^{-\left(R+S^{\prime}\right)}+\ldots+x^{R^{\prime}+S}\right)\left(x^{-T}+\ldots+1+\ldots+x^{T}\right)\right\}$
is the sum of all the coefficients of $f^{++} g$, and therefore

$$
\begin{equation*}
\mathbb{C}\left(f^{++g} h^{*}\right)=\sum_{r, s} a_{r} b_{s}=\sum_{r} a_{r} \sum_{s} b_{s}=\left(R+1+R^{\prime}\right)\left(S+1+S^{\prime}\right) \tag{10.7.3}
\end{equation*}
$$

The result follows from (10.7.2) and (10.7.3).
We now consider the general case in which

$$
R^{\prime}>0, \quad S^{\prime}>0, \quad \operatorname{Max}\left(R+S^{\prime}, \quad R^{\prime}+S\right)=n>T
$$

We assume that the result has been proved for

$$
\operatorname{Max}\left(R+S^{\prime}, \quad R^{\prime}+S\right)<n
$$

and argue inductively.
Let $x^{\rho}$ be the highest power in $f, x^{\sigma}$ the lowest in $g$, and write

$$
f-x^{\rho}=\phi, \quad g-x^{\sigma}=\psi, \quad f^{+}-x^{R^{\prime}}=\bar{\phi}, \quad+g-x^{-s^{\prime}}=\bar{\psi}
$$

Since $R^{\prime}>0, S^{\prime}>0$, none of these functions vanishes identically.

Then

$$
f g h=\left(\phi+x^{\rho}\right)\left(x^{\sigma}+\psi\right) h=\phi \psi h+\chi^{h},
$$

where

$$
\chi=x^{\sigma} \phi+x^{\rho+\sigma}+x^{\rho} \psi
$$

Since the highest power in $x^{\sigma} \phi$ is lower than $x^{\rho+\sigma}$, and the lowest power in $x^{\rho} \psi$ is higher than $x^{\rho+\sigma}$, there is no overlapping, and all coefficients in $\chi$ are 0 or 1 . Since the sum of the coefficients in $h$ is $2 T+1$, it follows that

$$
\begin{gather*}
\mathfrak{C}(\chi h) \leqq 2 T+1, \\
\mathfrak{C}(f g h) \leqq \mathfrak{C}(\phi \psi h)+2 T+1 . \tag{10.7.4}
\end{gather*}
$$

On the other hand
(10.7.5) $f^{++} g h^{*}=\left(\bar{\phi}+x^{R^{\prime}}\right)\left(x^{-s^{\prime}}+\psi\right) h^{*}=\bar{\phi} \bar{\psi} h^{*}+\bar{\chi} h^{*}$,
where

$$
\begin{aligned}
\bar{\chi} & =x^{-s^{\prime}} \bar{\phi}+x^{R^{\prime}-s^{\prime}}+x^{R^{\prime}} \bar{\psi} \\
& =x^{-R-s^{\prime}}+\ldots+x^{R^{\prime}-s^{\prime}-1}+x^{R^{\prime}-s^{\prime}}+x^{R^{\prime}-s^{\prime}+1}+\ldots+x^{R^{\prime}+s} .
\end{aligned}
$$

The sequence of exponents in $\bar{\chi}$ is an unbroken one, extending from $-R-S^{\prime}$ to $R^{\prime}+S$. We know that either $R+S^{\prime}$ or $R^{\prime}+S$ is greater than $T$. If $R+S^{\prime}>T$ then, by (10.6.1),

$$
R^{\prime}+S \geqq R+S^{\prime}-1 \geqq T,
$$

and so the unbroken sequence from $-T$ to $T$, of length $2 T+1$, is part of the sequence of exponents of $\bar{\chi}$. The same conclusion follows when $R^{\prime}+S>T$. Since $h^{*}$ has an unbroken sequence of exponents of length $2 T+1$, centred round the constant term, it follows that

$$
\mathfrak{C}\left(\bar{\chi} h^{*}\right)=2 T+1,
$$

and so, by (10.7.5), that

$$
\begin{equation*}
\mathfrak{C}\left(f^{++} g h^{*}\right)=\mathfrak{C}\left(\bar{\phi} \psi h^{*}\right)+2 T+1 . \tag{10.7.6}
\end{equation*}
$$

Now

$$
\phi^{+}(x)=x^{-\left(R^{\prime}-1\right)}+\ldots+x^{R}=\bar{\phi}\left(x^{-1}\right), \mathrm{a}
$$

Also

$$
{ }^{+} \psi(x)=x^{-s}+\ldots+x^{s^{\prime}-1}=\bar{\psi}\left(x^{-1}\right) .
$$

$$
\operatorname{Max}\left(R^{\prime}-1+S, R+S^{\prime}-1\right)=\operatorname{Max}\left(R^{\prime}+S, R+S^{\prime}\right)-1=n-1,
$$ and so, by our hypothesis,

$$
\begin{align*}
& \mathfrak{C}(\phi \psi h) \leqq \mathfrak{C}\left(\phi^{++} \psi h^{*}\right)=\mathfrak{C}\left\{\bar{\phi}\left(x^{-1}\right) \psi\left(x^{-1}\right) h^{*}(x)\right\}  \tag{10.7.7}\\
&=\mathfrak{C}\left\{\bar{\phi}(x) \Psi(x) h^{*}\left(x^{-1}\right)\right\}=\mathfrak{C}\left(\bar{\phi} \psi h^{*}\right) .
\end{align*}
$$

a It is not necessarily true that

$$
\phi^{+}(x)=\bar{\phi}(x)=x^{-R}+\ldots+x^{R^{\prime}-1},
$$

since, if $R^{\prime}=R$, this polynomial is 'overweighted' at the wrong end. When $R^{\prime}=R+1$, either formula is correct.

Finally, comparing (10.7.4), (10.7.6), and (10.7.7), we see that

$$
\mathfrak{C}(f g h) \leqq \mathfrak{C}\left(f^{++} g h^{*}\right),
$$

and the proof is completed.
10.8. Another proof of Theorem 371. There is another proof of Theorem 371 which is interesting in itself, although it cannot be extended to prove the more general Theorem 373.

We have to prove that, among the arrangements of the $x$ and $y$ which make $S$ a maximum, there is one in which
if

$$
\begin{align*}
x_{r}-x_{r^{\prime}} \geqq 0, & y_{s}-y_{s^{\prime}} \geqq 0  \tag{10.8.1}\\
\left|r^{\prime}\right|>|r|, & \left|s^{\prime}\right|>|s| \\
r^{\prime}=-r<0, & s^{\prime}=-s<0 .
\end{align*}
$$

or if
We may suppose on grounds of continuity that the $x, y$, and $c$ are all positive, that the $x$ and the $y$ are all different, and that the $c$ are different except in so far as they are restricted by $c_{-n}=c_{n}$, the condition of symmetry.

We shall denote an arrangement of the $x$ and $y$ generally by $A$. We say that $A$ is 'correct' if it satisfies (10.8.1); there is just one correct arrangement $C$. We say that $A$ is 'almost correct' if it satisfies (10.8.1) except perhaps when $r^{\prime}=-r$ or $s^{\prime}=-s$; there are, including $C, 2^{2 k}$ almost correct arrangements, and we denote the class of such arrangements by $C^{\prime}$. Finally we denote by $K$ the class of those $A$ which give the maximum value of $S$. We have to prove that $C$ is a $K$.

Given $p$, we can associate the $x$ and $y$ in pairs
(10.8.2) $\quad\left(x_{p-i}, x_{p+i}\right),\left(y_{p-j}, y_{p+j}\right) \quad(i, j=1,2,3, \ldots)$,
or in pairs
(10.8.3) $\quad\left(x_{p-i}, x_{p+i+1}\right), \quad\left(y_{p-j}, y_{p+j+1}\right) \quad(i, j=0,1,2, \ldots)$.

If a suffix falls outside the interval $(-k, k)$ then the corresponding $x$ or $y$ is to be replaced by 0 . In the first case the elements $x_{p}$ and $y_{p}$ are left unpaired, and the pairing may be made, by choice of appropriate $p, i, j$, to include any pair of elements the difference of whose ranks is positive and even. In the second case no elements are left unpaired, and the pairing may be made to include any pair the difference of whose ranks is odd. We use both pairings, and the arguments are essentially the same whichever is being used.

Consider, for example, the pairing (10.8.2), and suppose, to fix our ideas, that $p \geqq 0$, so that

$$
|p-i| \leqq|p+i|, \quad|p-j| \leqq|p+j| .
$$

We denote by $I, J$ those values of $i$ and $j$ for which

$$
x_{p-I}<x_{p+I}, \quad y_{p-J}<y_{p+J},
$$

so that the pairs corresponding to $I$ and $J$ do not satisfy (10.8.1). Such pairs we call 'wrong', others 'right'. If $p+i$ falls outside $(-k, k)$, but $p-i$ inside it, then $x_{p+i}$ is to be replaced by 0 , and the corresponding pair of $x$ is certainly right. Hence, except perhaps when $p=0$, there are $i$ which are not $I$ and $j$ which are not $J$.

If, for a given $p$ and a given pairing ${ }^{\text {a }}$, there are no wrong pairs, we say that $A$ is 'right with respect to $p$ ', and otherwise 'wrong with respect to $p^{\prime}$. It is clear that $C$ is right with respect to every $p$, and that any $C^{\prime}$ is right with respect to all $p$ except perhaps $p=0$ and the pairing (10.8.2). Further, any $A$ other than $C$ is wrong with respect to some $p$ and pairing, and any $A$ which is not a $C^{\prime}$ is wrong either with respect to a $p$ other than 0 or with respect to $p=0$ and the pairing (10.8.3).

We now (again envisaging the first pairing, and supposing $p \geqq 0$ ) consider the effect on $S$ of the substitution

$$
\Omega_{p} \quad\left(x_{p-I}, x_{p+I} ; y_{p-J}, y_{p+J}\right)
$$

which interchanges each pair $x_{p-I}, x_{p+I}$ and each pair $y_{p-J}, y_{p+J}$.
We divide $S$ into nine partial sums defined as follows:

```
\(S_{1}: r=p ; \quad s=p ;\)
\(S_{2}: \quad r=p ; \quad s=p-j, p+j(j \neq J) ;\)
\(S_{3}: \quad r=p-i, p+i(i \neq I) ; \quad s=p ;\)
\(S_{4}: \quad r=p ; \quad s=p-J, p+J ;\)
\(S_{5}: \quad r=p-I, p+I ; \quad s=p ;\)
\(S_{6}: \quad r=p-i, p+i(i \neq I) ; \quad s=p-j, p+j(j \neq J)\);
\(S_{7}: \quad r=p-i, p+i(i \neq I) ; \quad s=p-J, p+J ;\)
\(S_{8}: \quad r=p-I, p+I ; \quad s=p-j, p+j(j \neq J) ;\)
\(S_{9}: \quad r=p-I, p+I ; \quad s=p-J, p+J\).
```

It is plain, first, that $S_{1}, S_{2}, S_{3}$, and $S_{6}$ are not affected by $\Omega_{p}$.
${ }^{\text {a }}$ Either (10.8.2) or (10.8.3). In what follows 'right (or wrong) with respect to $p$ ' means always 'right (or wrong) with respect to $p$ and the pairing under consideration'.

Next

$$
S_{4}=x_{p} \sum_{J}\left(c_{J} y_{p-J}+c_{-J} y_{p+J}\right)
$$

is not affected because $c_{-J}=c_{J}$. Similarly $S_{5}$ and

$$
\begin{aligned}
S_{9}=\sum_{I, J}\left(c_{-I+J} x_{p-I} y_{p-J}+c_{-I-J}\right. & x_{p-I} y_{p+J} \\
& \left.+c_{I+J} x_{p+I} y_{p-J}+c_{I-J} x_{p+I} y_{p+J}\right)
\end{aligned}
$$

are not affected. It remains to consider $S_{7}$ and $S_{8}$.
The contribution of the pair $x_{p-i}, x_{p+i}$ to $S_{7}$ is

$$
x_{p-i} \sum_{J}\left(c_{-i+J} y_{p-J}+c_{-i-J} y_{p+J}\right)+x_{p+i} \sum_{J}\left(c_{i+J} y_{p-J}+c_{i-J} y_{p+J}\right)
$$

and the increment produced by $\Omega_{p}$ is

$$
-\left(x_{p-i}-x_{p+i}\right) \sum_{J}\left(c_{i-J}-c_{i+J}\right)\left(y_{p-J}-y_{p+J}\right)
$$

The total change in $S_{7}$ is the sum of this increment over $i \neq I$, and is positive, provided that there are $J$ and $i \neq I$, since the three differences written are respectively positive, positive, and negative. Hence $S_{7}$ is increased if there are $J$ and $i \neq I$; and similarly $S_{8}$ is increased if there are $I$ and $j \neq J$. Finally, $S$ is increased if either of these conditions is satisfied.

If $p \neq 0$, there are $i \neq I$ and $j \neq J$; and then $S$ is increased unless $A$ is right with respect to $p$. In any case, whatever $p, S$ is not diminished.

Suppose now that $A$ is not a $C^{\prime}$. Then $A$ is wrong with respect either to some $p \neq 0$ or to $p=0$ and the pairing (10.8.3). The argument above, or the similar argument based upon the pairing (10.8.3), then shows that $S$ is increased by $\Omega_{p}$ (or the corresponding substitution based upon the other pairing), and that $A$ is not a $K$. Hence the $K$ are included among the $C^{\prime}$. But if a $C^{\prime}$ is not $C$, then the substitution $\Omega_{0}$ replaces it by $C$ and does not diminish $S$; and therefore $C$ is a $K$. ${ }^{\text {a }}$

It does not seem to be possible to prove Theorem 373 by any equally simple argument based upon a substitution defined directly.
10.9. Rearrangements of any number of sets. There are

[^73]analogues of Theorem 373 for more than three sets ( $a$ ), ..., which may be deduced from Theorem 373 itself.
374. ${ }^{\mathrm{a}} \operatorname{If}(a),(b),(c),(d), \ldots$ are finite sets of non-negative numbers, and $(c),(d), \ldots$ are symmetrical, then
(10.9.1) $\sum_{r+s+t+u+\ldots=0}^{\Sigma} a_{r} b_{s} c_{t} d_{u} \ldots \leqq \sum_{r+s+t+u+\ldots=0}^{\Sigma} a_{r}{ }^{++} b_{s} c_{t}{ }^{*} d_{u}{ }^{*} \ldots$

We assume the theorem to be true when there are $k-1$ symmetrical sets (c), (d), ... involved, and prove that it is true when there are $k$. We shall make use of the following theorem, which is of some interest in itself.
375. If $\left(c^{*}\right),\left(d^{*}\right), \ldots$ are symmetrically decreasing sets, then the set ( $Q$ ) defined by

$$
\begin{equation*}
Q_{n}=\sum_{t+u+\ldots=n} c_{t} * d_{u} * \ldots \tag{1.9.2}
\end{equation*}
$$

is symmetrically decreasing.
It is enough to prove the theorem for two sets $\left(c^{*}\right),\left(d^{*}\right)$, since its truth in general then follows by repetition of the argument. We may agree that, when there is no indication to the contrary, sums involving several suffixes are extended over values of the suffixes whose sum vanishes.

It is plain that $Q_{-n}=Q_{n}$. Further, for any set ( $x$ ), we have

$$
\sum_{m} x_{m} Q_{m}=\Sigma x_{m} Q_{n}=\Sigma x_{m} c_{i}^{*} d_{u}^{*} \leqq \Sigma x_{m}^{+} c_{t}^{*} d_{u}^{*}=\sum_{m} x_{m}^{+} Q_{m},
$$

by Theorem 373. It follows, by Theorem 369, that the $Q_{m}$ are similarly ordered to the $x_{m}^{+}$, and therefore, since $Q_{m}$ is an even function of $m$, that the set is symmetrically decreasing.

This is the most elegant proof, but there is a simpler one which does not depend upon Theorem 373.

We drop the asterisks for convenience and suppose $n \geqq 0$. Then

$$
Q_{n}=\Sigma c_{n+r} d_{-r}+\Sigma c_{n+1-r} d_{r-1},
$$

the summations extending over $r \geqq 1$. Similarly

$$
Q_{n+1}=\Sigma c_{n+r} d_{1-r}+\Sigma c_{n+1-r} d_{r} .
$$

Subtracting, and using the equations $d_{-r}=d_{r}$ and $d_{1-r}=d_{r-1}$, we obtain

$$
\begin{aligned}
Q_{n}-Q_{n+1} & =\Sigma\left\{c_{n+r}\left(d_{-r}-d_{1-r}\right)+c_{n+1-r}\left(d_{r-1}-d_{r}\right)\right\} \\
& =\Sigma\left(c_{n+1-r}-c_{n+r}\right)\left(d_{r-1}-d_{r}\right) .
\end{aligned}
$$

Since $|n+1-r|<n+r$ for $n \geqq 0, r \geqq 1$, each term here is non-negative.

[^74]Returning to the proof of Theorem 374, we define $Q_{n}$ as in (10.9.2), and $P_{m}$ by

$$
P_{m}=\sum_{r+s=m} a_{r} b_{s}
$$

Then

$$
\Sigma a_{r} b_{s} c_{t} d_{u} \ldots=\Sigma P_{m} c_{t} d_{u} \ldots \leqq \Sigma P_{m}^{+} c_{t} * d_{u} * \ldots
$$ by the case $k-1$ of the theorem ${ }^{\text {a }}$. That is to say,

$$
\Sigma a_{r} b_{s} c_{t} d_{u} \ldots \leqq \Sigma P_{m}^{+} Q_{n}=\sum_{m} P_{m}^{+} Q_{m}=\sum_{m} P_{m} Q_{\phi(m)},
$$

where $\phi(m)$ is a permutation function for which $P_{m}=P_{\phi(m)}^{+}$; i.e.

$$
\begin{aligned}
& \Sigma a_{r} b_{s} c_{t} d_{u} e_{v} \ldots \leqq \Sigma a_{r} b_{s} Q_{\phi(m)} \leqq \Sigma a_{r}++b_{s} Q_{m}^{*} \\
&=\Sigma a_{r}^{++} b_{s} Q_{m}=\Sigma a_{r}++b_{s} c_{t}^{*} d_{u} * e_{v}^{*} \ldots,
\end{aligned}
$$

which is (10.9.1).
From Theorem 374 we can deduce ${ }^{b}$
376. Given any finite number of sets (a), (b), ..., we have

$$
\begin{aligned}
& \Sigma a_{r_{1}} a_{-r_{2}} b_{s_{1}} b_{-s_{2}} c_{t_{1}} c_{-t_{2}} \ldots \leqq \Sigma a_{r_{1}}^{+}+a_{r_{2}} b_{s_{1}}^{+}+b_{s_{2}} c_{t_{1}}^{+}+c_{t_{2}} \cdots \\
&=\Sigma a_{r_{1}}^{+} a_{-r_{2}}^{+} b_{s_{1}}^{+} b_{-s_{2}}^{+} c_{t_{1}}^{+} c_{-t_{2}}^{+} \cdots
\end{aligned}
$$

10.10. A further theorem on the rearrangement of any number of sets. In Theorems 373 and 374 two of the sets, (a) and (b), were arbitrary, but the remainder were subject to the condition of 'symmetry'. This restriction is essential; if (a), (b) and (c) are unrestricted, it is not possible to specify the maximal arrangement generally by means of the symbols $a^{+},+a, \ldots$ c

There is however a less precise theorem which is often equally effective in applications.
377. For any system of $k$ sets (a), (b), (c), ...

$$
\underset{r+s+t+\ldots=0}{\Sigma} a_{r} b_{s} c_{t} \ldots \leqq K(k) \underset{r+s+t+\ldots=0}{\Sigma} a_{r}^{+} b_{s}^{+} c_{t}^{+} \ldots,
$$

where $K=K(k)$ is a number depending only on $k$.
We suppose $k=3$; the argument is essentially the same in the general case.

We define the sets $\left(\beta^{*}\right),\left(\gamma^{*}\right)$ by

$$
\begin{array}{rr}
\beta_{m}^{*}=b_{m}^{+}, & \gamma_{m}^{*}=c_{m}^{+} \quad(m \geqq 0),  \tag{10.10.1}\\
\beta_{-m}^{*}=\beta_{m}^{*}, \quad \gamma_{-m}^{*}=\gamma_{m}^{*} \quad(m \geqq 0) ;
\end{array}
$$

a With $P, c, d, \ldots$ for $a, b, c, \ldots$ : since $(c)$ is symmetrical, $+c_{t}=c_{t}{ }^{*}$.
b Gabriel (3).
c See Theorem 388 at the end of the chapter.
and $(\beta)$ and $(\gamma)$ as the sets into which $\left(\beta^{*}\right)$ and $\left(\gamma^{*}\right)$ are changed by the permutations which change $\left(b^{+}\right)$into $(b)$ and $\left(c^{+}\right)$into (c) respectively. Then $(\beta)$ and $(\gamma)$ are symmetrical sets. Further, since $b_{-m}^{+} \leqq b_{m}^{+}$and $c_{-m}^{+} \leqq c_{m}^{+}$when $m \geqq 0$, we have

$$
b_{n}^{+} \leqq \beta_{n}^{*}, \quad c_{n}^{+} \leqq \gamma_{n}^{*}
$$

for all $n$, and so

$$
\begin{equation*}
b_{n} \leqq \beta_{n}, \quad c_{n} \leqq \gamma_{n} \tag{10.10.3}
\end{equation*}
$$

for all $n$.
We shall also require an inequality for $\beta_{m}^{*}$ and $\gamma_{m}^{*}$ with $m<0$. We have $b_{n}^{+} \leqq b_{-n+1}^{+}$and $c_{n}^{+} \leqq c_{-n+1}^{+}$for $n \geqq 1$, and so, by (10.10.1) and (10.10.2),
(10.10.4) $\quad \beta_{m}^{*} \leqq b_{m+1}^{+}, \quad \gamma_{m}^{*} \leqq c_{m+1}^{+} \quad(m<0)$.

Using (10.10.3) and the symmetry of $(\beta)$ and $(\gamma)$, we find

$$
S=\sum_{r+s+t=0} a_{r} b_{s} c_{t} \leqq \sum_{r+s+t=0} a_{r} \beta_{s} \gamma_{l} \leqq \sum_{r+s+t=0} a_{r}^{+} \beta_{s}^{*} \gamma_{t}^{*}
$$

by Theorem 373. The last sum is

$$
\left(\underset{s \geqq 0, t \geqq 0}{\Sigma}+\sum_{s<0, t \geqq 0}^{\Sigma}+\sum_{s \geqq 0, t<0}^{\Sigma}+\sum_{s<0, t<0}\right) a_{r}^{+} \beta_{s}^{*} \gamma_{t}^{*}
$$

and so, by (10.10.1) and (10.10.4),
(10.10.5)

$$
\begin{aligned}
S \leqq & \sum_{s \geqq 0, t \geqq 0} a_{r}^{+} b_{s}^{+} c_{t}^{+}+\sum_{s<0, t \geqq 0}^{\sum} a_{r}^{+} b_{s+1}^{+} c_{t}^{+} \\
& +\sum_{s \geqq 0, t<0} a_{r}^{+} b_{s}^{+} c_{t+1}^{+}+\sum_{s<0, t<0}^{\sum} a_{r}^{+} b_{s+1}^{+} c_{t+1}^{+}=S_{1}+S_{2}+S_{3}+S_{4}
\end{aligned}
$$

In $S_{2}, s<0$ and $r+s+t=0$, so that either $r>0$ or $t>0$. In the first case $a_{r}^{+} \leqq a_{r-1}^{+}$and in the second $c_{t}^{+} \leqq c_{t-1}^{+}$. Hence in any case, in $S_{2}$,
(10.10.6) $\quad a_{r}^{+} b_{s+1}^{+} c_{t}^{+} \leqq a_{r-1}^{+} b_{s+1}^{+} c_{t}^{+}+a_{r}^{+} b_{s+1}^{+} c_{t-1}^{+}$.

Similarly, in $S_{3}$,
(10.10.7) $\quad a_{r}^{+} b_{s}^{+} c_{t+1}^{+} \leqq a_{r-1}^{+} b_{s}^{+} c_{t+1}^{+}+a_{r}^{+} b_{s-1}^{+} c_{t+1}^{+}$.

Finally, in $S_{4}, s<0, t<0$ and $r+s+t=0$, so that $r \geqq 2$ and $a_{r}^{+} \leqq a_{r-2}^{+}$, and
(10.10.8)

$$
a_{r}^{+} b_{s+1}^{+} c_{t+1}^{+} \leqq a_{r-2}^{+} b_{s+1}^{+} c_{t+1}^{+} \cdot
$$

If now we substitute into (10.10.5) the upper bounds for the typical terms given by (10.10.6), (10.10.7), and (10.10.8), and
observe that, in these upper bounds, the sum of the suffixes is always 0 , we obtain

$$
S \leqq(1+2+2+1) \sum_{r+s+t=0} a_{r}^{+} b_{s}^{+} c_{t}^{+}=\underset{r+s+t=0}{6 \Sigma} a_{r}^{+} b_{s}^{+} c_{t}^{+}
$$

which proves the theorem.
10.11. Applications. These theorems have important applications to the theory of Fourier series. It is easy to deduce from Theorem $376^{\text {a }}$ that if

$$
f(\theta)=\sum_{-R}^{R} a_{r} e^{r \theta i}, \quad f^{+}(\theta)=\sum_{-R}^{R} \alpha_{r} e^{r \theta i},
$$

where $\alpha_{r}=\left|a_{r}\right|^{+}$, and $k$ is a positive integer, then

$$
\int_{-\pi}^{\pi}|f(\theta)|^{2 k} d \theta \leqq \int_{-\pi}^{\pi}\left|f^{+}(\theta)\right|^{2 k} d \theta ;
$$

and this relation between trigonometrical polynomials may be extended to functions represented by general Fourier series. Series of the type

$$
\Sigma \alpha_{r} e^{r \theta i}
$$

have particularly simple properties. They converge uniformly except at the origin and congruent points, where the function which they represent has in general an infinite peak; and the ratio

$$
\int_{-\pi}^{\pi}\left|f^{+}(\theta)\right|^{2 k} d \theta: \Sigma(|r|+1)^{2 k-2} \alpha_{r}^{2 k}
$$

lies between positive bounds depending only on $k$. We thus find, for example, that

$$
\int_{-\pi}^{\pi}|f(\theta)|^{2 k} d \theta \leqq K(k) \Sigma(|r|+1)^{2 k-2} \alpha_{r}^{2 k}
$$

For fuller developments see Hardy and Littlewood (9), Paley (3).
10.12. The rearrangement of a function. The theorems of $\S \S 10.1-10.10$ have analogues for functions of a continuous variable.

Suppose that $\phi(x)$ is non-negative and integrable in $(0,1)$, so that it is measurable and finite almost everywhere. If $M(y)$ is the measure of the set in which $\phi(x) \geqq y, M(y)$ is a decreasing function of $y$. The inverse $\bar{\phi}$ of $M$ is defined by

$$
\bar{\phi}\{M(y)\}=y ;
$$

and $\bar{\phi}(x)$ is a decreasing function of $x$ defined uniquely in $(0,1)$ except for at most an enumerable set of values of $x$, viz. those corre-

[^75]sponding to intervals of constancy of $M(y)$. We may complete the definition of $\bar{\phi}(x)$ by agreeing, for example, that
$$
\bar{\phi}(x)=\frac{1}{2}\{\bar{\phi}(x-0)+\bar{\phi}(x+0)\}
$$
at a point of discontinuity ${ }^{\text {a }}$.
We call $\bar{\phi}(x)$ the rearrangement of $\phi(x)$ in decreasing order. It is a decreasing function of $x$ which has, in general, an infinite peak at the origin.

The measure of the set in which $\bar{\phi}(x) \geqq y$ is $M(y) .^{\text {b }}$ It follows that the two (in general quite different) sets in which

$$
y_{1} \leqq \phi(x)<y_{2}, \quad y_{1} \leqq \Phi(x)<y_{2}
$$

have the same measure, and that the same is true of the sets in which

$$
\phi(x)>y, \quad \bar{\phi}(x)>y .
$$

We may say that the functions $\phi(x)$ and $\bar{\phi}(x)$ are 'equi-measurable'; they have equal integrals over $(0,1)$ and

$$
\int_{0}^{1} F(\bar{\phi}) d x=\int_{0}^{1} F(\phi) d x
$$

for any measurable $F$ for which the integrals exist.
We may define $\bar{\phi}(x)$ similarly for a $\phi(x)$ defined in any interval of $x$, provided that, if the interval is infinite, $M(y)$ is finite for every positive $y$.

If $\phi_{1}(x) \leqq \phi(x)$ then plainly $\bar{\phi}_{1}(x) \leqq \bar{\phi}(x)$. Suppose in particular that $\phi_{1}(x)$ is $\phi(x)$ in $E$ and zero in $C E$. Then

$$
\begin{equation*}
\int_{E} \phi(x) d x=\int \phi_{1}(x) d x=\int_{0}^{m E} \Phi_{1}(x) d x \leqq \int_{0}^{m E} \phi(x) d x . \tag{10.12.1}
\end{equation*}
$$

We shall use this inequality in $\S$ 10.19. In particular

$$
\begin{equation*}
\int_{0}^{x} \phi(t) d t \leqq \int_{0}^{x} \phi(t) d t \tag{10.12.2}
\end{equation*}
$$

if $\phi(x)$ is defined in $(0, a)$ and $0 \leqq x \leqq a$.
${ }^{\text {a }}$ Compare § 6.15.
${ }^{\mathrm{b}}$ This becomes obvious on drawing a figure. It must be remembered that $\bar{\phi}(x)$ may have intervals of constancy, corresponding to discontinuities of $M(y)$. It is however easy to prove that $M(y-0)=M(y)$ for all $y$, and so that the assertion in the text is true even for these exceptional $y$. In fact

$$
M\left(y-\frac{1}{n}\right)-M(y)=m S_{n}
$$

where $S_{n}$ is the set in which $y-n^{-1} \leqq \phi<y$, and the limit of $m S_{n}$ is zero.

Another type of rearrangement of a function will be important in what follows. Suppose, for example, that $\phi(x)$ is defined for all real, or almost all real $x$, and that $M(y)$ is finite for all positive $y$. We may define an even function $\phi^{*}(x)$ by agreeing that

$$
\phi^{*}\left\{\frac{1}{2} M(y)\right\}=y
$$

and that $\phi^{*}(-x)=\phi^{*}(x)$; or, what is the same thing, that $\phi^{*}(x)$ is even and

$$
\phi^{*}(x)=\bar{\phi}(2 x)
$$

for positive $x$. Then $\phi^{*}(x)$ decreases symmetrically on each side of the origin, where it has generally an infinite cuspidal peak. We call $\phi^{*}(x)$ the rearrangement of $\phi(x)$ in symmetrical decreasing order.
10.13. On the rearrangement of two functions. We begin by proving an integral inequality corresponding to Theorem 368.
378. Whether a is finite or infinite,

$$
\int_{0}^{a} \phi \psi d x \leqq \int_{0}^{a} \bar{\phi} \bar{\psi} d x
$$

We prove this by an argument similar to that of $\S 10.3$. In the first place, the theorem is true for functions which assume only the values 0 and 1. For suppose that $E$ and $F$ are the sets in which $\phi=1$ and $\psi=1$ respectively, and $\bar{E}, \bar{F}$ the analogous sets corresponding to $\bar{\phi}, \bar{\psi}$. Then the first integral is $m(\boldsymbol{E} F)$, the measure of the set $E F$ which is the product of $E$ and $F$, and

$$
m(E F) \leqq \operatorname{Min}(m E, m F)=\operatorname{Min}(m \bar{E}, m \bar{F})=m(\bar{E} \bar{F})
$$

Next, the theorem is true for functions which assume only a finite number of non-negative values. In fact, following the lines of $\S 10.3$, we can represent such a function $\phi$ in the form

$$
\phi=\alpha_{1} \phi_{1}+\alpha_{2} \phi_{2}+\ldots+\alpha_{n} \phi_{n}
$$

where the $\alpha$ are non-negative, the $\phi$ are always 0 or 1 , and

$$
\bar{\phi}=\alpha_{1} \bar{\phi}_{1}+\alpha_{2} \bar{\phi}_{2}+\ldots+\alpha_{n} \bar{\phi}_{n}
$$

The inequality then follows from a linear combination of inequalities already proved.

Finally, we prove the theorem in the general case by approximating to $\phi$ and $\psi$, by functions of the type just considered. We
do not give the last two stages of the proof in detail, since the arguments will recur in the proof of themore difficult Theorem 379.
10.14. On the rearrangement of three functions. We come now to what is our main object in these sections, the integral theorem corresponding to Theorems 372 and 373.
379. ${ }^{\text {a }}$ If $f(x), g(x)$, and $h(x)$ are non-negative, and $f^{*}(x), g^{*}(x)$, and $h^{*}(x)$ are the equi-measurable symmetrically decreasing functions, then

$$
\begin{align*}
I= & \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(x) g(y) h(-x-y) d x d y  \tag{10.14.1}\\
& \leqq \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f^{*}(x) g^{*}(y) h^{*}(-x-y) d x d y=I^{*}
\end{align*}
$$

We may plainly suppose that none of $f, g, h$ is nul. We may also replace $-x-y$ by $\pm x \pm y$ without changing the significance of the inequality.

We prove the inequality (1) for functions which are always 0 or $1,(2)$ for functions which take only a finite number of values, and (3) for general functions. As with Theorem 373, the whole difficulty lies in stage (1). We take this stage for granted for the moment and begin by showing that, if the theorem is true in this special case, it is true generally.

A function which takes only a finite number of non-negative values $0, a_{1}, a_{2}, \ldots, a_{n}$ can be expressed in the form

$$
f(x)=\alpha_{1} f_{1}(x)+\alpha_{2} f_{2}(x)+\ldots+\alpha_{n} f_{n}(x)
$$

where the $\alpha$ are positive, the $f_{j}$ take only the values 0 and 1 , and

$$
f_{1} \geqq f_{2} \geqq \ldots \geqq f_{n} .
$$

For we may suppose $0<a_{1}<a_{2}<\ldots<a_{n}$, take

$$
\alpha_{1}=a_{1}, \quad \alpha_{2}=a_{2}-a_{1}, \quad \ldots, \quad \alpha_{n}=a_{n}-a_{n-1}
$$

and

$$
\begin{aligned}
& f_{1}=1 \quad\left(f \geqq a_{1}\right), \quad=0 \quad\left(f<a_{1}\right), \\
& f_{2}=1 \quad\left(f \geqq a_{2}\right), \quad=0 \quad\left(f<a_{2}\right),
\end{aligned}
$$

A moment's consideration shows that we then have also

$$
\begin{gathered}
f *(x)=\alpha_{1} f_{1}^{*}(x)+\alpha_{2} f_{2}^{*}(x)+\ldots+\alpha_{n} f_{n}^{*}(x) . \\
\text { a F. Riesz (8). }
\end{gathered}
$$

If we suppose that each of $f, g, h$ takes only a finite number of values, and decompose them in this way, then (10.14.1) follows from the combination of similar inequalities involving triads $f_{i}, g_{j}, h_{k}{ }^{\text {a }}{ }^{\text {a }}$

To pass from this case to the general case, we approximate to $f, g, h$ by functions which take only a finite number of values. We can approximate to $f$, for example, by the function $f_{n}$ defined by

$$
f_{n}=\frac{k}{n}\left(\frac{k}{n} \leqq f<\frac{k+1}{n}, k=0,1,2, \ldots, n^{2}-1\right), \quad f_{n}=n(f \geqq n) ;
$$

and to $g$ and $h$ similarly. Then $f_{n} \leqq f, f_{n}^{*} \leqq f^{*}$, and similarly for $g$ and $h$. Hence (assuming that the theorem has been proved for the special type of functions) we have

$$
I_{n}=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f_{n}(x) g_{n}(y) h_{n}(-x-y) d x d y \leqq I_{n} * \leqq I^{*}
$$

and so $I=\lim I_{n} \leqq I^{*}$.
It remains to prove the theorem in the special case when $f, g, h$ assume only the values 0 and 1 . It is however convenient first to make a further reduction of the problem.

First, we may suppose that the sets $F, G$, and $H$ in which $f, g$, and $h$ assume the value 1 are finite. If two of these sets are infinite, then two of $f^{*}, g^{*}$, and $h^{*}$ are 1 for all $x$, in which case $I^{*}=\infty^{\text {b }}$ and there is nothing to prove. Suppose then that just one of the sets, say $F$, is infinite. Let $F_{N}$ be the part of $F$ in $(-N, N)$, let $N$ be the smallest number for which $m F_{N} \geqq 2 n$, and define $f_{n}$ as being $f$ in $F_{N}$ and 0 outside. Then (assuming the theorem to have been proved when the sets are finite)

$$
\begin{aligned}
& I_{n}=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f_{n} g h d x d y \leqq \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f_{n}^{*} g^{*} h^{*} d x d y \\
= & \int_{-n}^{n} \int_{-\infty}^{\infty} f^{*} g^{*} h^{*} d x d y \leqq \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f^{*} g^{*} h^{*} d x d y=I^{*},
\end{aligned}
$$

and so

$$
I=\lim I_{n} \leqq I^{*}
$$

Suppose then that $f(x)$ assumes the value 1 in a set $E$ of finite measure. We can represent $E$ in the form $\varepsilon+e-e^{\prime}$, where $\varepsilon$ is a

[^76]finite set of non-overlapping intervals, and $e$ and $e^{\prime}$ are sets of arbitrarily small measure ${ }^{\text {a }}$; and the sets in which $g$ and $h$ assume the value 1 can be represented similarly. It is also plain that, since $f, g$, and $h$ do not exceed 1 , small sets $e, \ldots$ make a small difference in the integrals $I$ and $I^{*}$. We may therefore suppose that the sets in which $f=1, g=1$, and $h=1$ are finite sets of intervals; if the theorem has been proved in this case, its truth in the more general case follows by approximation.

Next we may suppose, on similar grounds, that the ends of all the intervals are rational; and then, by a change of variable, that they are integral. The theorem is thus reduced to dependence upon the case in which each of the sets in which $f=1, g=1$, or $h=1$ consists of a finite number of intervals ( $m, m+1$ ), where $m$ is an integer.

Finally we may suppose, if we please, that the number of intervals in any or all of the sets is even, since we can replace each interval by two by bisecting it and effecting another change of variable.
10.15. Completion of the proof of Theorem 379. It is convenient to replace $f(x)$ by $f(-x)$, as plainly we may without affecting the result. If we do this, write $s, t$ for $x, y$ and then make the substitution $s=x-t$, we obtain

$$
\begin{equation*}
I=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(t-x) g(t) h(-x) d x d t=\int_{-\infty}^{\infty} h(-x) \chi(x) d x \tag{10.15.1}
\end{equation*}
$$

$I^{*}=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f^{*}(t-x) g^{*}(t) h^{*}(-x) d x d t=\int_{-\infty}^{\infty} h^{*}(-x)^{*} \chi(x) d x$, where

$$
\begin{equation*}
\chi(x)=\int_{-\infty}^{\infty} f(t-x) g(t) d t, \quad{ }^{*} \chi(x)=\int_{-\infty}^{\infty} f^{*}(t-x) g^{*}(t) d t .{ }^{\mathrm{b}} \tag{10.15.3}
\end{equation*}
$$

We suppose for the moment merely that $f, g, h$ are characteristic

[^77]functions of sets (functions assuming the values 0 and 1 only), without using the further simplifications shown to be permissible at the end of $\S 10.14$. We denote the sets in which $f, \ldots, f^{*}, \ldots$ assume the value 1 by $F, \ldots, F^{*}, \ldots$; each function vanishes outside the corresponding set, and $F^{*}, \ldots$ are intervals symmetric about the origin. We suppose that
$$
m F=m F^{*}=2 R, \quad m G=m G^{*}=2 S, \quad m H=m H^{*}=2 T
$$

With this notation ${ }^{\text {a }}$, we have (10.15.4)

$$
I \leqq \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(t-x) g(t) d x d t=\int_{-\infty}^{\infty} f(-s) d s \int_{-\infty}^{\infty} g(t) d t=4 R S
$$

$$
\begin{equation*}
I^{*}=\int_{-T}^{T} * \chi(x) d x \tag{10.15.5}
\end{equation*}
$$

If $x$ is fixed, and $t-x$ describes the set $F$, then $t$ describes a set $F_{x}$ obtained by translating $F$ through a distance $x$. If we define $F^{*}{ }_{x}=\left(F^{*}\right)_{x}$ similarly, then the functions (10.15.3) may be written in the form

$$
\begin{equation*}
\chi(x)=m\left(F_{x} G\right), \quad * \chi(x)=m\left(F_{x}^{*} G^{*}\right) \tag{10.15.6}
\end{equation*}
$$

From this formula we can calculate ${ }^{*} \chi(x)$. Let us suppose, as we may, that

$$
\begin{equation*}
R \leqq S \tag{10.15.7}
\end{equation*}
$$

Then * $\chi(x)$ is continuous,

$$
(10.15 .8) * \chi(x)=0 \quad(|x| \geqq R+S), \quad * \chi(x)=2 R \quad(|x| \leqq S-R)
$$

and ${ }^{*} \chi(x)$ is linear in the intervals $(-R-S,-S+R)$ and ( $S-R, R+S$ ). The graph of ${ }^{*} \chi(x)$ is shown in Fig. 3.

Suppose now that

$$
\begin{equation*}
R+S \leqq T \tag{10.15.9}
\end{equation*}
$$

Then it follows from (10.15.8) that

$$
I^{*}=\int_{-T}^{T} * \chi(x) d x=\int_{-R-S}^{R+S} * \chi(x) d x=4 R S
$$

and the result of the theorem follows from (10.15.4). We have thus proved the theorem under the restriction (10.15.9). It is also plainly true if $R=0$ or $S=0$ (when $F$ or $G$ is nul).
a Chosen to emphasize the parallelism of the argument with that of $\S \S 10.6-7$.

So far $F, G, H$ have been arbitrary sets of finite measure. We now make the further specialisation explained at the end of $\S 10.14$, supposing that $F, G, H$ are sets of intervals ( $m, m+1$ ), the numbers of intervals being $2 R, 2 S, 2 T$ respectively. We may if we please suppose these numberseven, but we shall argue inductively, and it is more convenient to adopt a slightly more general hypothesis and to suppose only that $2 R+2 S+2 T$ is even. In these circumstances $R, S$, and $T$ are not necessarily integral, but $2 R$, $2 S, 2 T$ and

$$
\begin{equation*}
\mu=R+S-T=R+S+T-2 T \tag{10.15.10}
\end{equation*}
$$



Fig. 3. Graph of ${ }^{*} \chi$ and ${ }^{*} \chi_{1}$.
are integral. We have already proved that the theorem is true if $\mu \leqq 0$, and it is also true if $R=0$ or $S=0$. It is therefore enough to establish its truth when

$$
\begin{equation*}
\mu=n>0, \quad R>0, \quad S>0, \tag{10.15.11}
\end{equation*}
$$

on the assumption that it is true when $\mu=n-1$.
We denote by $F_{1}$ the set derived from $F$ by omitting the last interval of $F$ on the right; similarly $G_{1}$ is $G$ less the last interval on the right of $G$. Generally, sets, functions, or numbers with suffix 1 are derived from $F_{1}$ and $G_{1}$ as the corresponding sets, functions, or numbers without suffixes are derived from $F$ and $G$; thus $f_{1} *$ is the rearrangement of $f_{1}$, the characteristic function of $F_{1}$, and ${ }^{*} \chi_{1}(x)$ is the 'Faltung' of $f_{1}{ }^{*}$ and $g_{1}{ }^{*} . F_{1}{ }^{*}$ is the interval

$$
\left(-R+\frac{1}{2}, \quad R-\frac{1}{2}\right),
$$

and generally, $R$ and $S$ are replaced by $R-\frac{1}{2}$ and $S-\frac{1}{2}$ when we pass from $F, G$ to $F_{1}, G_{1}$. By the inductive hypothesis

$$
\begin{equation*}
I_{1} \leqq I_{1}{ }^{*} . \tag{10.15.12}
\end{equation*}
$$

The function ${ }^{*} \chi_{1}(x)$ vanishes for $|x| \geqq R+S-1$, is equal to $2 R-1$ for $|x| \leqq S-R$, and is linear in the intervals remaining; and $T \leqq R+S-1$, by (10.15.10) and (10.15.11). Hence ${ }^{\text {a }}$

$$
{ }^{*} \chi(x)-{ }^{*} \chi_{1}(x)=1
$$

for $-T \leqq x \leqq T$, and, by (10.15.5),

$$
\begin{equation*}
I^{*}-I_{1}^{*}=\int_{-T}^{T}\left\{* \chi(t)-* \chi_{1}(t)\right\} d t=2 T \tag{10.15.13}
\end{equation*}
$$

We have now to consider

$$
\begin{equation*}
I-I_{1}=\int_{-\infty}^{\infty} h(-x)\left\{\chi(x)-\chi_{1}(x)\right\} d x . \tag{10.15.14}
\end{equation*}
$$

Here, after (10.15.6),

$$
\begin{equation*}
\chi(x)-\chi_{1}(x)=m\left(F_{x} G\right)-m\left(F_{1 x} G_{1}\right) . \tag{10.15.15}
\end{equation*}
$$

This function is plainly linear in any interval ( $m, m+1$ ), and therefore assumes its extrema for integral values of $x$. Suppose then that $x$ is integral. In this case the set $F_{x} G$ is composed of full intervals ( $m, m+1$ ), and, when we remove the intervals of $F_{x}$ and $G$ furthest to the right, either one or no interval of $F_{x} G$ is lost, one if the extreme interval of either set coincides with an interval of the other set, and none if there is no such coincidence ${ }^{b}$. Hence $\chi(x)-\chi_{1}(x)$ is 1 or 0 for integral $x$, and therefore

$$
\begin{equation*}
0 \leqq \chi(x)-\chi_{1}(x) \leqq 1 \tag{10.15.16}
\end{equation*}
$$

for all $x$.
From (10.15.14) and (10.15.16) it follows that

$$
\begin{equation*}
0 \leqq I-I_{1}=\int_{H}\left\{\chi(-x)-\chi_{1}(-x)\right\} d x \leqq \int_{H} d x=2 T, \tag{10.15.17}
\end{equation*}
$$

and from (10.15.13) and (10.15.17) that

$$
\begin{equation*}
I-I_{1} \leqq I^{*}-I_{1}^{*} . \tag{10.15.18}
\end{equation*}
$$

Finally (10.15.12) and (10.15.18) give $I \leqq I^{*}$; and this completes the proof ${ }^{c}$.

[^78]10.16. An alternative proof. The proof of Theorem 379 given by Riesz is also very interesting. We can simplify it by reducing the theorem, as in $\S 10.14$, to the case in which $f, g, h$ are each equal to $l$ in a finite set of intervals and to 0 elsewhere. We represent the variables $x, y, z$ on the sides of an equilateral triangle, taking the middle point of each side as the origin and the positive directions on the sides cyclically. Then $x+y+z=0$ is the condition that the points $x, y, z$ on the sides should be the three orthogonal projections of a point of the plane ${ }^{\text {a }}$.

The functions $f(x), g(y), h(z)$ are the characteristic functions of three sets $E_{1}, E_{2}, E_{3}$, each consisting of a finite number of non-overlapping intervals, and $f^{*}(x), g^{*}(y), h^{*}(z)$ are the characteristic functions of the three intervals $E_{1}{ }^{*}, E_{2}{ }^{*}, E_{3}{ }^{*}$ of lengths $E_{1}, E_{2}, E_{3}, \mathrm{~b}$ symmetrically disposed about the three origins. If $E_{123}$ is the set of those points of the plane whose three projections belong to $E_{1}, E_{2}$, and $E_{3}$, and $E^{*}{ }_{123}$ is defined similarly, thenc

$$
I=\sin \frac{1}{3} \pi E_{123}, \quad I^{*}=\sin \frac{1}{3} \pi E^{*}{ }_{123},
$$

and what we have to prove is that

$$
\begin{equation*}
E_{123} \leqq E^{*}{ }_{123} . \tag{10.16.1}
\end{equation*}
$$

The figure $E^{*}{ }_{123}$ is defined by drawing six lines perpendicular to the sides, and is a hexagon unless one of $E_{1}, E_{2}, E_{3}$ is greater than the sum of the other two, in which case it reduces to a parallelogram. We begin by proving (10.16.1) in the latter case. Suppose for example that $E_{3} \geqq E_{1}+E_{2}$. Then $E^{*}{ }_{123}$ reduces to
the extreme intervals on the right. It would not have been true if we had removed two arbitrary intervals. Suppose, for example, that each of $F$ and $G$ is the interval $(-4,4)$, that $F_{4}$ consists of the two intervals $(-4,-2)$ and $(2,4)$, and $G_{4}$ of the interval ( $-2,2$ ). We can pass from $F, G$ to $F_{4}, G_{4}$ in four steps, taking away one unit interval from each set at each step; but

$$
\chi(0)-\chi_{4}(0)=8,
$$

instead of being less than or equal to 4 . The same example shows that Zygmund's assertion (1, 176) 'those [the values] of $\phi(x)$ in $(-\infty, \infty)$ increase at most by 2 ' is untrue unless his construction is restricted in a way which he does not state explicitly. It is essential to go closely into detail at this point, since it is the kernel of the proof.
${ }^{\text {a }}$ If $P$ is the point in question and $G$ is the centre of the triangle, then

$$
x+y+z=P G\left\{\cos \alpha+\cos \left(\alpha+\frac{2}{3} \pi\right)+\cos \left(\alpha+\frac{4}{3} \pi\right)\right\}=0 .
$$

${ }^{\mathrm{b}}$ We use $E_{1}$ both for the set $E_{1}$ and for its measure.
c $E_{123}$, when used as a measure, is of course a plane measure.
$E^{*}{ }_{12}$, the set of points projecting into $E_{1}{ }^{*}$ and $E_{2}{ }^{*}$ on two of the sides, while $E_{123}$ is included in the set $E_{12}$ defined similarly. Hence

$$
E_{123} \leqq E_{12}=\operatorname{cosec} \frac{1}{3} \pi \cdot E_{1} E_{2}=\operatorname{cosec} \frac{1}{3} \pi \cdot E_{1} * E_{2}^{*}=E_{12}^{*}=E^{*}{ }_{123}
$$

This proves the theorem when $E^{*}{ }_{123}$ is a parallelogram.
Passing to the case of the hexagon, suppose for example that

$$
E_{3}>E_{1} \geqq E_{2}, \quad E_{3}<E_{1}+E_{2}
$$

We define sets

$$
E_{1}(t), \quad E_{2}(t), \quad E_{3}(t)
$$

and corresponding intervals $E_{1}{ }^{*}(t), E_{2}{ }^{*}(t), E_{3}{ }^{*}(t)$, by subtracting from each $E_{j}$ a set of measure $t$ at each end a. If $t$ increases from 0 to

$$
t_{0}=\frac{1}{2}\left(E_{1}+E_{2}-E_{3}\right)
$$

$E_{1}(t), E_{2}(t), E_{3}(t)$ decrease from $E_{1}, E_{2}, E_{3}$ to sets $E_{1}\left(t_{0}\right), E_{2}\left(t_{0}\right)$, $E_{3}\left(t_{0}\right)$ whose measures satisfy

$$
E_{1}\left(t_{0}\right)+E_{2}\left(t_{0}\right)=E_{3}\left(t_{0}\right)
$$

The hexagon then reduces to a parallelogram, so that

$$
\begin{equation*}
E_{123}\left(t_{0}\right) \leqq E_{123}{ }_{123}\left(t_{0}\right) \tag{10.16.2}
\end{equation*}
$$

If we can prove also that

$$
\begin{equation*}
E_{123}-E_{123}\left(t_{0}\right) \leqq E^{*}{ }_{123}-E^{*}{ }_{123}\left(t_{0}\right) \tag{10.16.3}
\end{equation*}
$$

our conclusion will follow by addition.
We prove (10.16.3) by comparing the derivatives of

$$
\phi(t)=-E_{123}(t), \quad \phi^{*}(t)=-E^{*}{ }_{123}(t) .
$$

In the first place, the difference between $E^{*}{ }_{123}(t)$ and $E^{*}{ }_{123}(t+h)$ is a hexagonal ring whose area is $h P(t)+O\left(h^{2}\right)$, where $P(t)$ is the perimeter of the hexagon corresponding to the value $t$, and so

$$
\frac{d \phi^{*}}{d t}=P(t)=\operatorname{cosec} \frac{1}{3} \pi\left\{E_{1}(t)+E_{2}(t)+E_{3}(t)\right\}
$$

On the other hand the three sets

$$
E_{1}(t)-E_{1}(t+h), \quad E_{2}(t)-E_{2}(t+h), \quad E_{3}(t)-E_{3}(t+h)
$$

a That is to say

$$
E_{1}=E_{1}^{\prime}(t)+E_{1}(t)+E_{1}^{\prime \prime}(t),
$$

where $E_{1}{ }^{\prime}(t)$ lies to the left, and $E_{1}{ }^{\prime \prime}(t)$ to the right, of $E_{1}(t)$, and

$$
m E_{1}^{\prime}(t)=m E_{1}^{\prime \prime}(t)=t .
$$

consist in all of six intervals, each of length $h$, for small $h$. The twelve perpendiculars to the sides of the triangle drawn through the ends of the six intervals define a hexagonal ringa which includes the whole of $E_{123}(t)-E_{123}(t+h)$. The derivative $\phi^{\prime}(t)$ is the total length of those parts of the outer boundary of this ring which also belong to $E_{123}(t)$. Projecting these parts of the


Fig. 4. The decrement of $E_{123}(t)$.
boundary on to the sides of the triangle, as indicated in the figure, we see that

$$
\frac{d \phi}{d t} \leqq \operatorname{cosec} \frac{1}{3} \pi\left\{E_{1}(t)+E_{2}(t)+E_{3}(t)\right\}=\frac{d \phi^{*}}{d t}
$$

From this (10.16.3) follows by integration, and this completes the proof of the theorem.

[^79]10.17. Applications. The special case of Theorem 379 corresponding to Theorem 371 is
380. If $h(x)$ is symmetrically decreasing, then
\[

$$
\begin{aligned}
& I=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f(x) g(y) h(x-y) d x d y \\
& \leqq \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} f^{*}(x) g^{*}(y) h(x-y) d x d y=I^{*}
\end{aligned}
$$
\]

We shall now apply Theorems 371 and 380 to the special cases
and

$$
c_{r-s}=|r-s|^{-\lambda}
$$

$$
h(x-y)=|x-y|^{-\lambda} .
$$

381. If

$$
a_{r} \geqq 0, \quad b_{s} \geqq 0,
$$

(10.17.1) $p>1, \quad q>1, \quad \frac{1}{p}+\frac{1}{q}>1, \quad \lambda=2-\frac{1}{p}-\frac{1}{q}$
(so that $0<\lambda<1$ ), and

$$
\begin{gathered}
\Sigma a_{r}^{p}=A, \quad \Sigma b_{s}^{q}=B \\
T=\Sigma \Sigma^{\prime} \frac{a_{r} b_{s}}{|r-s|^{\lambda}} \leqq K A^{1 / p} B^{1 / q}
\end{gathered}
$$

where the dash implies that $r \neq s$, and $K=K(p, q)$ depends on $p$ and $q$ only.
382. If $f(x) \geqq 0, g(y) \geqq 0, p$ and $q$ satisfy (10.17.1), and

$$
\begin{gathered}
\int_{-\infty}^{\infty} f^{p}(x) d x=F, \quad \int_{-\infty}^{\infty} g^{q}(y) d y=G \\
I=\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \frac{f(x) g(y)}{|x-y|^{\lambda}} d x d y \leqq K F^{1 / p} G^{1 / q}
\end{gathered}
$$

then
The proofs of the two theorems are practically the same. We give that of Theorem 382. ${ }^{\text {a }}$

It is plain, after Theorem 380, that we may replace $f$ and $g$ by $f^{*}$ and $g^{*}$. We then divide $I$ into four parts corresponding to the four quadrants of integration. The north-east and south-west parts are equal, and so are the north-west and south-east parts, and the two latter do not exceed the two former ${ }^{\text {b }}$. We need
${ }^{\text {a }}$ For that of Theorem 381 see Hardy, Littlewood, and Pólya (1); for a deduction of Theorem 382 from Theorem 381 see Hardy and Littlewood (6).
${ }^{b}$ The north-west and south-east parts could be accounted for by the easier argument of $\S 9.14$.
therefore only consider the north-east part. Hence, changing our notation again, it is sufficient to prove that

$$
J=\int_{0}^{\infty} \int_{0}^{\infty} \frac{f(x) g(y)}{|x-y|^{\lambda}} d x d y \leqq K F^{1 / p} G^{1 / q}
$$

where now $f$ and $g$ are positive and decreasing, and $F$ and $G$ are defined by integrals over $(0, \infty)$. We write

$$
J=J_{1}+J_{2}
$$

where $J_{1}$ and $J_{2}$ are integrals over the octants $y \leqq x$ and $x \leqq y$ respectively.

We have

$$
J_{1}=\int_{0}^{\infty} f(x) d x \int_{0}^{x} \frac{g(y)}{(x-y)^{\lambda}} d y
$$

Since $g(y)$ decreases, and $(x-y)^{-\lambda}$ increases, in $(0, x)$,

$$
x \int_{0}^{x} \frac{g(y)}{(x-y)^{\lambda}} d y \leqq \int_{0}^{x} g(y) d y \int_{0}^{x} \frac{d y}{(x-y)^{\lambda}}=\frac{x^{1-\lambda}}{1-\lambda} g_{1}(x),
$$

say, by Theorem 236. Hence

$$
J_{1} \leqq \frac{1}{1-\lambda} \int_{0}^{\infty} f(x) g_{1}(x) x^{-\lambda} d x
$$

By Theorem 189

$$
\begin{equation*}
J_{1} \leqq \frac{1}{1-\lambda} F^{1 / p}\left(\int_{0}^{\infty} g_{1}^{p^{\prime}}(x) x^{-p^{\prime} \lambda} d x\right)^{1 / p^{\prime}} \tag{10.17.2}
\end{equation*}
$$

But $p^{\prime}>q$, by (10.17.1), and

$$
g_{1}(x)=\int_{0}^{x} g(y) d y \leqq G^{1 / q} x^{1 / q^{\prime}}
$$

again by Theorem 189. Hence
(10.17.3)

$$
g_{1}^{p^{\prime}}(x) x^{-p^{\prime} \lambda} \leqq g_{1}^{q}(x)\left(G^{1 / q} x^{1 / q^{\prime}}\right)^{p^{\prime}-q} x^{-p^{\prime} \lambda}=G^{\left(p^{\prime}-q\right) / q}\left\{\frac{g_{1}(x)}{x}\right\}^{q}
$$

(since

$$
\left.\frac{p^{\prime}-q}{q^{\prime}}-\lambda p^{\prime}=\frac{p^{\prime}-q}{q^{\prime}}-\left(\frac{1}{p^{\prime}}+\frac{1}{q^{\prime}}\right) p^{\prime}=-\frac{q}{q^{\prime}}-\mathbf{1}=-q\right)
$$

From (10.17.2) and (10.17.3) it follows that

$$
J_{1} \leqq K F^{1 / p} G^{\left(p^{\prime}-q\right) / p^{\prime} q}\left\{\int_{0}^{\infty}\left(\frac{g_{1}(x)}{x}\right)^{q} d x\right\}^{1 / p^{\prime}} \leqq K F^{1 / p} G^{1 / q}
$$

by Theorem 327.
The discussion of $J_{2}$ is similar, and the theorem follows.
383. Suppose that $f(x)$ is non-negative and $L^{p}$, where $p>1$, in $(0, \infty)$, that

$$
\begin{equation*}
0<\alpha<\frac{1}{p}, \quad q=\frac{p}{1-\alpha p} \tag{10.17.4}
\end{equation*}
$$

and that

$$
\begin{equation*}
f_{\alpha}(x)=\frac{1}{\Gamma(\alpha)} \int_{0}^{x} f(y)(x-y)^{\alpha-1} d y \tag{10.17.5}
\end{equation*}
$$

Then $f_{\alpha}(x)$ is $L^{q}$ in $(0, \infty)$ and

$$
\begin{equation*}
\int_{0}^{\infty} f_{\alpha}^{q} d x \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{q / p} \tag{10.17.6}
\end{equation*}
$$

where

$$
K=K(p, \alpha)=K(p, q)
$$

Suppose that $g(x)$ is any function of $L^{q^{\prime}}$, and that

$$
\lambda=1-\alpha=1-\frac{1}{p}+\frac{1}{q}=2-\frac{1}{p}-\frac{1}{q^{\prime}}
$$

By Theorem 382,

$$
\int_{0}^{\infty} \int_{0}^{\infty} \frac{g(x) f(y)}{|x-y|^{\lambda}} d x d y \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{q^{\prime}} d x\right)^{1 / q^{\prime}}
$$

and a fortiori

$$
\begin{gathered}
\int_{0}^{\infty} f_{\alpha}(x) g(x) d x=\frac{1}{\Gamma(\alpha)} \int_{0}^{\infty} g(x) d x \int_{0}^{x} \frac{f(y)}{(x-y)^{\lambda}} d y \\
\leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{q^{\prime}} d x\right)^{1 / q^{\prime}}
\end{gathered}
$$

Since this is true for all $g$, it follows, by Theorem 191, that

$$
\left(\int_{0}^{\infty} f_{\alpha}^{q} d x\right)^{1 / q} \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}
$$

which is (10.17.6).
The proof shows that the result is also true when $f_{\alpha}(x)$ is defined by

$$
f_{\alpha}(x)=\frac{1}{\Gamma(\alpha)} \int_{x}^{\infty} f(y)(y-x)^{\alpha-1} d y
$$

Theorem 383 embodies a result in the theory of 'fractional integration'. Liouville (1) and Riemann (1, 331-344) defined the integral $f_{\alpha}(x)$ of $f(x)$, of order $\alpha$, as

$$
\begin{equation*}
f_{\alpha}(x)=\frac{1}{\Gamma(\alpha)} \int_{a}^{x} f(y)(x-y)^{\alpha-1} d y \tag{10.17.7}
\end{equation*}
$$

The lower limit $a$ is the 'origin of integration'; a change of origin changes $f_{\alpha}$ in a manner which is not trivial formally, though unimportant for
theorems of the type considered here. It is easily deduced from Theorem $383^{a}$ that, if $f$ is $L^{p}$ in $(a, b)$, where $-\infty<a<b \leqq \infty, \alpha<1 / p$, and $f_{\alpha}$ is the integral of $f$, of order $\alpha$ and with origin $a$, then $f_{\alpha}$ is $L^{q}$ in $(a, b)$. When $\alpha>1 / p$, $f_{\alpha}$ is continuous, and indeed belongs to the 'Lipschitz class' of order $\alpha-1 / p$.

In applications of the theory, $f$ is usually periodic. It was observed by Weyl (3) that the reference to an origin $a$ is in this case inappropriate: Weyl accordingly modified the definition as follows. If we suppose that the mean value of $f$ over a period is zero (a condition which we can always satisfy by subtracting an appropriate constant from $f$ ), then

$$
\int_{-\infty}^{x} f(y)(x-y)^{a-1} d y
$$

converges at the lower limit, and we may take $a=-\infty$ in (10.17.7). Our theorem concerning the Lebesgue classes may be extended to this case also.
10.18. Another theorem concerning the rearrangement of a function in decreasing order. The theorem with which we end is important primarily for its function-theoretic applications, but the proof which we give ${ }^{b}$ is interesting independently.

The theorem may be stated in two forms.
384. Suppose that $f(x)$ is non-negative and integrable in a finite interval $(0, a)$, that $\bar{f}(x)$ is the rearrangement of $f(x)$ in decreasing order, that

$$
\begin{equation*}
\Theta(x)=\Theta(x, f)=\operatorname{Max}_{0 \leq \xi<x} \frac{1}{x-\xi} \int_{\xi}^{x} f(t) d t, \tag{10.18.1}
\end{equation*}
$$

and that $\bar{\Theta}(x)$ is the rearrangement of $\Theta(x)$ in decreasing order. Then

$$
\begin{equation*}
\bar{\Theta}(x) \leqq \frac{1}{x} \int_{0}^{x} \bar{f}(t) d t \tag{10.18.2}
\end{equation*}
$$

for $0<x \leqq a$.
385. Suppose that $f(x)$ satisfies the conditions of Theorem 384, and that $s(y)$ is any increasing function of $y$ defined for $y \geqq 0$. Then

$$
\begin{equation*}
\int_{0}^{a} s\{\Theta(x)\} d x \leqq \int_{0}^{a} s\left\{\frac{\mathbf{1}}{x} \int_{0}^{x} \bar{f}(t) d t\right\} d x . \tag{10.18.3}
\end{equation*}
$$

We begin with two preliminary remarks.

[^80](1) We shall prove Theorem 384 first and deduce Theorem 385. Since $\Theta(x)$ and $\bar{\Theta}(x)$ are equimeasurable,
$$
\int_{0}^{a} s\{\Theta(x)\} d x=\int_{0}^{a} s\{\bar{\Theta}(x)\} d x .
$$

Hence (10.18.3) follows from (10.18.2).
That (10.18.2) follows from (10.18.3), so that the two forms of the theorem are equivalent, is a little less obvious, but is proved in Theorem 392. ${ }^{\text {a }}$ The first implication is sufficient for our purpose here, since it is in the second form that the theorem is used in the applications.
(2) If

$$
\Theta_{0}(x)=\Theta_{0}(x, f)=\frac{1}{x} \int_{0}^{x} f(t) d t
$$

then

$$
\Theta_{0}(x, f) \leqq \Theta_{0}(x, \bar{f})=\Theta(x, \bar{f})=\frac{1}{x} \int_{0}^{x} \bar{f}(t) d t,
$$

by (10.12.2), and

$$
\begin{equation*}
\int_{0}^{a} s\left\{\Theta_{0}(x)\right\} d x \leqq \int_{0}^{a} s\left\{\frac{1}{x} \int_{0}^{x} \bar{f}(t) d t\right\} d x . \tag{10.18.4}
\end{equation*}
$$

This, a much more trivial inequality than (10.18.3), is the analogue for integrals of Theorem 333.
10.19. Proof of Theorem 384. We suppose, as we may, that $a=1$.

We consider a point $x_{0}$ for which
write

$$
\begin{equation*}
\bar{\Theta}\left(x_{0}\right)=p+\epsilon \quad(p>0, \epsilon>0), \tag{10.19.1}
\end{equation*}
$$

and consider the set $E$ defined by

$$
\begin{equation*}
0 \leqq x \leqq 1, \quad \Theta(x)>p . \tag{10.19.2}
\end{equation*}
$$

Since $\Theta(x)$ and $\bar{\Theta}(x)$ are equimeasurable, $E$ has the same measure as the set in which $\bar{\Theta}(x)>p$. This set is at least as large as the set

[^81]in which $\bar{\Theta}(x) \geqq p+\epsilon$, and the measure of this last set is, after (10.19.1), at least $x_{0}$. Hence
\[

$$
\begin{equation*}
x_{0} \leqq m E \tag{10.19.3}
\end{equation*}
$$

\]

Now the set $E$ is composed of those points $x$ for which

$$
\begin{equation*}
\frac{1}{x-\xi} \int_{\xi}^{x} f(t) d t>p \tag{10.19.4}
\end{equation*}
$$

for some $\xi=\xi(x)<x$. We can write (10.19.4) in the form

$$
\begin{equation*}
\int_{0}^{x} f(t) d t-p x>\int_{0}^{\xi} f(t) d t-p \xi \tag{10.19.5}
\end{equation*}
$$

or
(10.19.6)

$$
g(x)>g(\xi)
$$

say. Thus $E$ is the set of points in which a certain continuous function $g(x)$ assumes a value greater than some at any rate of the values which it has assumed before. This property enables us to characterise the structure of $E$.

Lemma A. The set $E$ is composed of a finite or enumerable system of non-overlapping intervals $\left(\alpha_{k}, \beta_{k}\right)$. All of these intervals are open, and

$$
g\left(\alpha_{k}\right)=g\left(\beta_{k}\right)
$$

except possibly when $x=1$ is a point of $E$, in which case there is one interval $\left(\alpha_{k}, 1\right)$ closed on the right, and $g\left(\alpha_{k}\right) \leqq g(1)$, though $g\left(\alpha_{k}\right)$ is not necessarily equal to $g(1)$. ${ }^{\text {a }}$

In the first place, since $g(x)$ is continuous, $E$ is an open set (except possibly for the point $x=1$ ). Hence $E$ is a set of intervals $\left(\alpha_{k}, \beta_{k}\right)$, open if $\beta_{k}<1$.

If $\beta_{k}<1$ then $\beta_{k}$ is not a point of $E$, and

$$
\begin{equation*}
g\left(\alpha_{k}\right) \geqq g\left(\beta_{k}\right) \tag{10.19.7}
\end{equation*}
$$

by the definition of $E$.
Next, suppose that $\alpha_{k}<x_{1}<\beta_{k}$, and consider the minimum of $g(x)$ in the interval $0 \leqq x \leqq x_{1}$. This minimum cannot be attained for $\alpha_{k}<x \leqq x_{1}$, since all such $x$ belong to $E$, and so $g(x)>g(\xi)$ for some $\xi<x$. Hence it is attained for $x \leqq \alpha_{k}$. But $\alpha_{k}$ is not a point of $E$, and therefore $g\left(\alpha_{k}\right) \leqq g(x)$ for all these $x$. Hence the mini-

[^82]mum is attained at $\alpha_{k}$, and $g\left(\alpha_{k}\right) \leqq g\left(x_{1}\right)$. Making $x_{1} \rightarrow \beta_{k}$, we obtain
\[

$$
\begin{equation*}
g\left(\alpha_{k}\right) \leqq g\left(\beta_{k}\right) ; \tag{10.19.8}
\end{equation*}
$$

\]

and this, with (10.19.7), proves the lemma ${ }^{a}$.
We can now prove Theorem 384. We may write (10.19.8) in the form

$$
p\left(\beta_{k}-\alpha_{k}\right) \leqq \int_{\alpha_{k}}^{\beta_{k}} f(x) d x
$$

and from this it follows that

$$
p . m E=p \Sigma\left(\beta_{k}-\alpha_{k}\right) \leqq \Sigma \int_{\alpha_{k}}^{\beta_{k}} f(x) d x=\int_{E} f(x) d x .
$$

Hence, by (10.12.1),

$$
\begin{equation*}
p \cdot m E \leqq \int_{E} f(x) d x \leqq \int_{0}^{m E} \bar{f}(x) d x ; \tag{10.19.9}
\end{equation*}
$$

and hence, by (10.19.1),

$$
\begin{equation*}
\bar{\Theta}\left(x_{0}\right)-\epsilon=p \leqq \frac{1}{m E} \int_{0}^{m E} \bar{f}(x) d x . \tag{10.19.10}
\end{equation*}
$$

Finally, since $\bar{f}(x)$ decreases, it follows from (10.19.10) and (10.19.3) that

$$
\bar{\Theta}\left(x_{0}\right)-\epsilon \leqq \frac{1}{x_{0}} \int_{0}^{x_{0}} \bar{f}(x) d x .
$$

Since $\epsilon$ is arbitrary, this gives (10.18.1), with $x_{0}$ for $x$.
The function-theoretic applications of Theorems 384 and 385 arise as follows. Suppose that $f(\theta)$ is integrable and has the period $2 \pi$, that

$$
M(\theta)=M(\theta, f)=\operatorname{Max}_{0<|t| \leqq \pi} \frac{1}{t} \int_{0}^{t} f(\theta+u) d u
$$

and that $N(\theta)$ is the similar function formed with $|f(\theta+u)|$. These functions are of the same type as the $\Theta(x)$ of Theorem 384, but are generated by means taken to either side of $\theta$.

Consider now the integral

$$
\begin{equation*}
h(\theta, p)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} f(\theta+t)_{\chi}(t, p) d t \tag{i}
\end{equation*}
$$

where $\chi$ is a kernel which involves a parameter $p$ and satisfies the conditions
(ii)

$$
\chi(t, p) \geqq 0, \quad \frac{1}{2 \pi} \int_{-\pi}^{\pi} \chi(t, p) d t=1
$$

The standard examples of such kernels are the 'Poisson kernel'

$$
\chi=\frac{1-r^{2}}{1-2 r \cos t+r^{2}},
$$

${ }^{\text {a }}$ The argument here is due to M. Riesz (see F. Riesz, 10).
in which $p=r$ is positive and less than 1 , and the 'Fejér kernel'

$$
\chi=\frac{\sin ^{2} \frac{1}{2} n t}{n \sin ^{2} \frac{1}{2} t},
$$

in which $p=n$ is a positive integer. The corresponding values of $h$ are $u(r, \theta)$, the harmonic function defined by the 'Poisson integral' of $f(\theta)$, and $\sigma_{n}(\theta)$, the Cesàro mean, of order 1, of the Fourier series of $f(\theta)$.

Suppose now (a) that $f(\theta)$ belongs to $L^{k}$, where $k>1$, and (b) that $\chi$ satisfies the additional condition

$$
\begin{equation*}
\frac{1}{2 \pi} \int_{-\pi}^{\pi}\left|t \frac{\partial \chi}{\partial t}\right| d t \leqq A \tag{iii}
\end{equation*}
$$

where $A$ is independent of $p$. It follows from Theorems 385, with $s(y)=y^{k}$, and 327, that $M(\theta)$ also belongs to $L^{k} .{ }^{\text {a }}$ And it is easy to deduce from (i), (ii), and (iii) that

$$
|h(\theta, p)| \leqq A M(\theta)
$$

where $A$ is again independent of $p$. Hence $h$ has a majorant (independent of $p$ ) of the class $L^{k}$.

It is easily verified that the Poisson kernel satisfies (iii). Hence $u(r, \theta)$ possesses a majorant $U(\theta)$ of the class $L^{k}$. The same is true of $\sigma_{n}(\theta)$, but in this case the proof is not quite so simple, since the Fejer kernel does not satisfy (iii). We can however prove that $\left|\sigma_{n}(\theta)\right| \leqq A N(\theta)$, and similar conclusions follow. All this is set out in detail by Hardy and Littlewood (8).

## MISCELLANEOUS THEOREMS AND EXAMPLES

386. If $c_{2} \geqq c_{3} \geqq \ldots \geqq c_{2 n} \geqq 0$ and the sets ( $a$ ), (b) are non-negative and given except in arrangement, then

$$
\sum_{r=1}^{n} \sum_{s=1}^{n} c_{r+s} a_{r} b_{s}
$$

is a maximum when (a) and (b) are both in decreasing order.
[F. Wiener (1).]
387. It is not true that

$$
\sum_{r+s+t=0} a_{r} b_{s} c_{t} \leqq \sum_{r+s+t=0} a_{r}^{+} b_{s}^{+} c_{t}^{+} .
$$

[Trivial: take $(a),(b),(c)$ to be $(0,2,1),(1,2,0),(1,2,1)$. Then

$$
\left.\Sigma a_{r} b_{s} c_{t}=14, \quad \Sigma a_{r}+b_{s}+c_{t}^{+}=12 .\right]
$$

388. There are sets $(a),(b),(c)$ such that none of the eight sums

$$
\Sigma a^{+} b^{+} c^{+}, \quad \Sigma+a b^{+} c^{+}, \quad \Sigma a^{++} b c^{+}, \quad \ldots, \quad \Sigma+a^{+} b^{+} c
$$

gives the maximal sum $\Sigma a b c$.
[Suppose $0<h<1$ and $\epsilon$ positive and sufficiently small; and take (a) to be $0,0,0,1,2$, (b) to be $h-\epsilon, h, h+\epsilon, 1,1$, and (c) to be formed of any five different elements.]
a See Theorem 398 below.
389. If

$$
M(x)=\Sigma|r| x_{r}, \quad M(y)=\Sigma|s| y_{s},
$$

and $p \neq 0$, then the substitution $\Omega_{p}$ of $\S 10.8$ decreases $\mu=M(x)+M(y)$.
[The theorem is trivial, but may be used to construct another proof of Theorem 371, which follows the general lines of that in $\S 10.8$ but is free from an appeal to 'continuity'.

We use $A, C, C^{\prime}, K$ as in $\S 10.8$; there may now be more than one arrangement $C$. We define $L$ as the sub-class of $K$ formed by those members of $K$ for which $\mu$ is least. If $p \neq 0$, and $A$ is wrong with respect to $p$, then $\Omega_{p}$ decreases $\mu$ and does not diminish $S$. Hence any $A$ of $L$ is a $C^{\prime}$; and we can then show as in § 10.8 that $L$ includes a C.]
390. In the notation of Theorem 373
for every $n$.

$$
\underset{r+s+t=n}{\Sigma} a_{r} b_{s} c_{t} \leqq{\underset{r+s+t=0}{\Sigma} a_{r}^{+}+b_{s} c_{t} *}^{*}
$$

[Corollary of Theorem 373.]
391. If $(a),\left(a^{\prime}\right),(b),\left(b^{\prime}\right),(c)$, and $\left(c^{\prime}\right)$ are six sets of positive numbers subject to (10.5.1), then

$$
\underset{r+s+t=0}{\Sigma} a_{r} a_{r}{ }^{\prime} b_{s} b_{s}{ }^{\prime} c_{t} c_{t}^{\prime} \leqq \sum_{r+s+t=0} a_{r} * a_{r}{ }^{\prime *} b_{s} * b_{s}{ }^{\prime} * c_{t} * c_{t}{ }^{\prime *}
$$

[Corollary of Theorem 372 if first reduced, by the method of §10.3, to the special case in which every number is 0 or l.]
392. If $f$ and $g$ are non-negative, and

$$
\begin{equation*}
\int_{0}^{a} s\{f(x)\} d x \leqq \int_{0}^{a} s\{g(x)\} d x \tag{i}
\end{equation*}
$$

for every positive and increasing $s(y)$, then
(ii)

$$
\bar{f} \leqq \bar{g}
$$

except perhaps for an enumerable set of values of $x$.
[This is the theorem referred to in §10.18, as proving the equivalence of Theorems 384 and 385. It is an analogue of Theorem 107.

Since the integrals (i) are unaltered when we replace $f$ and $g$ by $\bar{f}$ and $\bar{g}$, we may suppose $f$ and $g$ themselves decreasing, so that $f=\bar{f}, g=\bar{g}$ (except perhaps in an enumerable set of points).

If (ii) is not true for almost all $x$, we can find a $b$ and a $c$ such that

$$
\begin{equation*}
b<c, \quad f(c)>g(b) . \tag{iii}
\end{equation*}
$$

For, if this were not so, we should have $f(b+0) \leqq g(b)$ for all $b$, and $f(b) \leqq g(b)$ at all points of continuity of the functions, and therefore except in an enumerable set.

Supposing then that $b$ and $c$ satisfy (iii), we choose $r$ so that

$$
g(b)<r<f(c)
$$

and define $s(y)$ by

$$
s(y)=0(y<r), \quad s(y)=1 \quad(y \geqq r) .
$$

Then $\quad \int_{0}^{a} s\{f(x)\} d x=\int_{f \geqq r} d x \geqq c>b \geqq \int_{g \geqq r} d x=\int_{0}^{a} s\{g(x)\} d x$,
in contradiction with (i).]
393. If $a_{1}, a_{2}, \ldots, a_{N}$ are non-negative,

$$
\Theta(n)=\Theta(n, a)=\operatorname{Max}_{1 \leqq \nu \leqq n} \frac{a_{\nu}+a_{\nu+1}+\ldots+a_{n}}{n-\nu+1},
$$

and a bar denotes a rearrangement in decreasing order (a notation opposite to that of § 10.1), then

$$
\Theta(n) \leqq \frac{\bar{a}_{1}+\bar{a}_{2}+\ldots+\bar{a}_{n}}{n} \quad(1 \leqq n \leqq N) .
$$

394. If the conditions of Theorem 393 are satisfied, and $s(y)$ is a positive increasing function of $y$, then

$$
\sum_{1}^{N} s\{\Theta(n)\} \leqq \sum_{1}^{N} s\left(\frac{\bar{a}_{1}+\bar{a}_{2}+\ldots+\bar{a}_{n}}{n}\right) .
$$

[The last two theorems are the analogues for finite sums of Theorems 384 and 385, and the reader will find it instructive to prove them by an adaptation of the argument of $\mathfrak{\xi} \S 10.18-19$. The earlier proofs of Hardy and Littlewood and of Gabriel are referred to in §10.18.]
395. If $\quad c_{1} \geqq c_{2} \geqq \ldots \geqq c_{p}>0, \quad d_{1} \geqq d_{2} \geqq \ldots \geqq d_{a}>0$;
$e_{1}, e_{2}, \ldots, e_{p+a}$ is the aggregate of the $c$ and $d$ rearranged in decreasing order;

$$
C_{n}=c_{1}+c_{2}+\ldots+c_{n},
$$

and $D_{n}$ and $E_{n}$ aredefinedsimilarly; and $s(y)$ is positive andincreasing; then

$$
\begin{gathered}
s\left(C_{1}\right)+s\left(\frac{C_{2}}{2}\right)+\ldots+s\left(\frac{C_{p}}{p}\right)+s\left(D_{1}\right)+s\left(\frac{D_{2}}{2}\right)+\ldots+s\left(\frac{D_{q}}{q}\right) \\
\leqq s\left(E_{1}\right)+s\left(\frac{E_{2}}{2}\right)+\ldots+s\left(\frac{E_{p+q}}{p+q}\right) .
\end{gathered}
$$

[This is a special case of Theorem 394. For a direct proof by induction, due to Chaundy, see Hardy and Littlewood (8): the theorem is one of the lemmas on which they based their proof of Theorem 394.]
396. If $p, q, P, Q$ are positive integers and $s(y)$ is positive and increasing, then

$$
\begin{gather*}
\sum_{1}^{p} s\left(\frac{p}{n}\right)+\sum_{1}^{q} s\left(\frac{q}{n}\right) \leqq \sum_{1}^{p+q} s\left(\frac{p+q}{n}\right),  \tag{i}\\
p+P \\
\sum_{p+1}^{p} s\left(\frac{p}{n}\right)+\sum_{q+1}^{q+Q} s\left(\frac{q}{n}\right) \leqq{ }^{p+q+P+Q} \sum_{p+q+1} s\left(\frac{p+q}{n}\right), \\
\sum_{1}^{\infty} s\left(\frac{p}{n}\right)+\sum_{1}^{\infty} s\left(\frac{q}{n}\right) \leqq \sum_{1}^{\infty} s\left(\frac{p+q}{n}\right) .
\end{gather*}
$$

[(i) and (ii) follow from Theorem 395 by appropriate specialisation, and (iii), which is true whether $p$ and $q$ are integers or not, is a corollary. A case of (iii) is

$$
\frac{x^{1 / a}}{1-x^{1 / a}}+\frac{x^{1 / b}}{1-x^{1 / b}} \leqq \frac{x^{1 /(a+b)}}{1-x^{1 /(a+b)}} \quad(a>0, b>0,0<x<1):
$$

this may naturally be proved independently (and with ' <'), for example as an application of Theorem 103.]
397. If $a, b, \alpha, \beta$ are positive and $s$ positive and increasing, then

$$
\int_{a}^{a+a} s\left(\frac{a}{x}\right) d x+\int_{b}^{b+\beta} s\left(\frac{b}{x}\right) d x \leqq \int_{a+b}^{a+b+a+\beta} s\left(\frac{a+b}{x}\right) d x
$$

398. If $k>1$, and $\Theta(x)$ is defined as in Theorem 384, then

$$
\int_{0}^{a} \Theta^{k}(x) d x \leqq\left(\frac{k}{k-1}\right)^{k} \int_{0}^{a} f^{k}(x) d x
$$

[From Theorems 385 and 327. There is of course a corresponding theorem for finite sums. This theorem has particularly important applications.]
399. In order that an integrable function $\phi(x)$ should have the property

$$
\int_{0}^{1} s(x) \phi(x) d x \geqq 0
$$

for all positive, increasing, and bounded $s(x)$, it is necessary and sufficient that

$$
\int_{x}^{1} \phi(t) d t \geqq 0 \quad(0 \leqq x \leqq 1) .
$$

[To prove the condition necessary, specialise $s(x)$ appropriately; to prove it sufficient, integrate partially or use the second mean value theorem. The condition is certainly satisfied if there is a $\xi$ between 0 and 1 such that $\phi(x) \geqq 0$ for $x>\xi, \phi(x) \leqq 0$ for $x<\xi$, and

$$
\int_{0}^{1} \phi(x) d x=0 .
$$

Theorem 397 is a special case of this theorem (after a simple transformation).]
400. If $E$ and $\xi$ are functions of $x$ subject to
then

$$
\begin{gathered}
\int_{0}^{1}\left\{\frac{E(x)-E(\xi)}{x-\xi}\right\}^{k} d x \leqq \frac{k E(1)-E^{k}(1)}{k-1} \quad(k>1), \\
\int_{0}^{1} \frac{E(x)-E(\xi)}{x-\xi} d x \leqq E(1)\left\{1+\log \frac{1}{E(1)}\right\} .
\end{gathered}
$$

[Suppose that $f(x)$ is always 0 or 1 , and that $E(x)$ is the measure of the part of $(0, x)$ in which $f(x)=1$, and apply Theorem 385.]
401. If $\quad p>1, q>1, \frac{1}{p}+\frac{1}{q} \geqq 1, \lambda=2-\frac{1}{p}-\frac{1}{q}$,

$$
h<1-\frac{1}{p}, \quad k<1-\frac{1}{q}, \quad h+k \geqq 0,
$$

and $h+k>0$ if $\frac{1}{p}+\frac{1}{q}=1$, then

$$
\int_{0}^{\infty} \int_{0}^{\infty} \frac{f(x) g(y)}{x^{h} y^{k}|x-y|^{\lambda-h-k}} d x d y \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{1 / p}\left(\int_{0}^{\infty} g^{q} d x\right)^{1 / q}
$$

[Here, and in Theorems 402 and 403, $K$ denotes a positive number depending only on the parameters of the theorem (here $p, q, h, k)$.]
402. If

$$
p>1, \quad 0 \leqq \alpha<\frac{1}{p}, \quad p \leqq q \leqq \frac{p}{1-\alpha p}
$$

then

$$
\int_{0}^{\infty} x^{-(p-q+p q a) / p} f_{a}^{q} d x \leqq K\left(\int_{0}^{\infty} f^{p} d x\right)^{q / p},
$$

where $f_{\alpha}$ is defined as in (10.17.5). The result is still true if $\alpha \geqq 1 / p$, when the second condition on $q$ may be omitted.
[For the last two theorems see Hardy and Littlewood (6). The case $q=p$ gives

$$
\int_{0}^{\infty}\left(x^{-a} f_{\alpha}\right)^{p} d x \leqq K \int_{0}^{\infty} f^{p} d x:
$$

compare Theorem 329.]
403. The result of Theorem 383 is not necessarily true when $p=1$.
[Define $f(x)$ by

$$
f(x)=\frac{1}{x}\left(\log _{\frac{1}{x}}\right)^{-\beta}\left(0<x \leqq \frac{1}{2}\right), \quad=0\left(x>\frac{1}{2}\right)
$$

where $\beta>1$. Then

$$
\begin{aligned}
& \qquad f_{\alpha}(x)=K \int_{0}^{x} \frac{1}{y}\left(\log \frac{1}{y}\right)^{-\beta}(x-y)^{\alpha-1} d y \\
& >K x^{\alpha-1} \int_{0}^{x} \frac{1}{y}\left(\log \frac{1}{y}\right)^{-\beta} d y=K x^{\alpha-1}\left(\log \frac{1}{x}\right)^{1-\beta} . \\
& \quad p=1, \quad q=\frac{1}{1-\alpha} ;
\end{aligned}
$$

$$
\left.\frac{\beta-1}{1-\alpha}>1, \quad \beta>2-\alpha .\right]
$$

404. Suppose that $f(x)$ is defined in $(-1,1)$ and has a continuous derivative $f^{\prime}(x)$ which vanishes only at a finite number of points, and that

$$
f(x) \geqq 0, \quad f(-1)=f(1)=0 .
$$

Then the length of the curve $y=f(x)$ is greater than that of $y=f^{*}(x)$, unless $f(x)=f^{*}(x)$.
[See Steiner (1, iI, 265). If $0<y<Y=\operatorname{Max} f$ then (except perhaps for a finite number of values of $y$ ) the equation $y=f(x)$ has an even number $2 n$ (depending on $y$ ) of roots. If we denote these roots, in ascending order, by $x_{1}, x_{2}, \ldots, x_{2 n}$, and the derivative of $x_{\nu}$ with respect to $y$ by $x_{\nu}{ }^{\prime}$, then, by Theorem 25,

$$
2 \int_{0}^{Y}\left\{1+\left[\frac{1}{2} \Sigma(-1)^{\nu} x_{\nu}^{\prime}\right]^{2}\right\}^{\frac{1}{2}} d y \leqq \int_{0}^{Y} \Sigma\left(1+x_{\nu}^{\prime 2}\right)^{\frac{1}{2}} d y
$$

There is equality only if $n=1$ for all $y$, and $x_{1}=-x_{2}$.]
405. Suppose that $f(x, y) \geqq 0$ for all $x, y$, and that the measure $M(z)$ of the set in which $f(x, y) \geqq z$ is finite for all positive $z$; define $\rho(z)$ by

$$
M(z)=\pi \rho^{2}
$$

and write

$$
f^{*}(x, y)=\rho^{-1}\left\{\sqrt{ }\left(x^{2}+y^{2}\right)\right\},
$$

where $\rho^{-1}$ is the inverse of $\rho$. Then (under appropriate conditions of regularity) the area of the surface $z=f(x, y)$ is greater than that of

$$
z=f^{*}(x, y) .
$$

[See Schwarz (1). The theorem is important in itself and interesting because it involves a two-dimensional analogue of the notion of $f^{*}(x)$.]

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[^0]:    a There are exceptions, as for example in §§ 8.8-8.17. There the 'positive' cases of the theorems discussed are relatively trivial.

[^1]:    a See Artin and Schreier (1).

[^2]:    ${ }^{\text {a }}$ Compare § 2.19.

[^3]:    a Except perhaps some parts of Ch. IV.

[^4]:    a If we admitted infinite values, there would be a corresponding case for positive $r$, viz. $r>0$, some $a$ infinite, $\mathfrak{M}_{r}=\infty$.

[^5]:    a This is what is usually called Cauchy's inequality: see Cauchy (1, 373). The corresponding inequality for integrals (Theorem 181) is usually called Schwarz's inequality, though it seems to have been stated first by Buniakowsky: see Buniakowsky (1, 4), Schwarz (2, 251).

[^6]:    a Hölder (1). Hölder states the theorem in a less symmetrical form given a little earlier by Rogers (1).

[^7]:    a Pringsheim (1), Jensen (2). Pringsheim attributes his proof to Lüroth.
    b Compare the remarks on this proof in §1.4.

[^8]:    a Hadamard (1) considers determinants with complex constituents. Theorem 30 was found earlier by Kelvin and proved by Muir (1).
    b See Bôcher (1, 171).

[^9]:    ${ }^{\text {a }}$ See also A. L. Dixon (1).
    ${ }^{\text {b }}$ Oppenheim (2). e Minkowski (2).
    ${ }^{a}$ Fischer (1). e See Bôcher (1, 171).

[^10]:    ${ }^{\text {a }}$ Thus $\Sigma d_{i k} x_{i} x_{k}$ is formed from $\Sigma c_{i k} x_{i} x_{k}$ by replacing $x_{i}, x_{k}(i, k=1,2, \ldots, r)$ by $-x_{i},-x_{k}$ (and is therefore positive if $\sum c_{i k} x_{i} x_{k}$ is positive).
    ${ }^{\mathrm{b}}$ See Wedderburn (1).

[^11]:    a Theorem 45 is due substantially to Muirhead (2); but Muirhead considers only integral $\alpha$.

[^12]:    a $\alpha_{l}-\alpha_{l}{ }^{\prime}$ is the first negative difference, $\alpha_{k}-\alpha_{k}^{\prime}$ the last positive difference which precedes it. The text assumes $l-k>1$; the case $l-k=1$ is easier.
    b Possibly both.

[^13]:    a This proof was known before Muirhead's work; see Hurwitz (1).
    b Communicated to us by Prof. I. Schur.

[^14]:    a Hilbert (1). b See, for example, Bôcher (1, 144-154).

[^15]:    ${ }^{a}$ See Theorem $58 . \quad$ b Hausdorff (1). Hausdorff has $n=2$.
    c Some of the theorems which follow here are mere exercises for the reader, but most have some independent interest.

[^16]:    ${ }^{\text {a }}$ In this chapter we define $\mathfrak{M}_{\phi}$ directly, and deduce its properties from the definition. In Ch. VI (§§ 6.19-6.22) we shall show how $\mathfrak{M}_{\phi}$ may be defined 'axiomatically', that is to say by prescription of its characteristic properties.
    b Knopp (2), Jessen (2).

[^17]:    a Provided $x \neq 1, y \neq 1$. Since (3.3.7) is plainly true when $x$ or $y$ is 1 , the exception is irrelevant.
    b Cauchy (1, 103-105).

[^18]:    a Here we follow the lines of $\S 2.6$ (ii). For a proof following Cauchy's argument more directly, see Jensen (2).

[^19]:    a M. Riesz (1), Jessen (2).
    ${ }^{\text {b }}$ See, for example, § 8.13.

[^20]:    a Jessen (2, 3).
    b We have actually proved more in regard to the necessity of the condition: see our remarks on Theorem 83.

[^21]:    ${ }^{\text {a }}$ The important case in practice is that in which (as stated in the theorem) $\phi^{\prime \prime}$ exists in the open interval. We usually wish, however, to assert convexity in the closed interval. Since $\phi^{\prime \prime} \geqq 0, \phi^{\prime}$ and $\phi$ are monotonic near the ends of the interval and tend to finite or infinite limits; $\phi^{\prime}$ may tend to $-\infty$ at the left-hand end and to $+\infty$ at the right-hand end, and $\phi$ may tend to $+\infty$ at either. The function will be convex in the closed interval if its value at each end is not less than its limit at that end.

[^22]:    a It would be sufficient to consider rectangular domains, but convexity is the natural limitation to impose on $D$. It is not part of our programme to consider questions of analysis situs connected with convex or general domains.
    b There is a wider generalisation of the notion of a convex function of a single variable which is important in the theory of functions but with which we shall not be concerned. The function $\Phi(x, y)$ is subharmonic if its value at the centre of any circle does not exceed its average over the circumference. In particular $\Phi$ is subharmonic if it is twice differentiable and

    $$
    \nabla^{2} \Phi=\Phi_{x x}+\Phi_{y y} \geqq 0
    $$

    For the theory of subharmonic functions see F. Riesz (5, 9), Montel (1).

[^23]:    a Jensen (2): Jensen does not refer to the necessity of the condition.

[^24]:    a We take for granted the obvious extensions of § 3.12 from two to $n$ variables.

[^25]:    ${ }^{\text {a }}$ The first results of this character are due to Bosanquet (1): Bosanquet considers case (II).

[^26]:    ${ }^{\text {a }}$ Schur (2) proves that (2) is a sufficient condition, and the remark concerning the case of equality is also due to him. For the complete theorem, see Hardy, Littlewood, and Pólya (2). Karamata (1) considers condition (1).

[^27]:    a See also Hahn (1,581).

[^28]:    a For a formal proof, use Theorem 183.

[^29]:    ${ }^{\text {a }}$ Maclaurin (2). See Pólya and Szegö (1, II, 45-47 and 230-232).
    b That is to say, all equations whose coefficients are not all zero.
    c Newton (1, 173). For further references, see § 2.22.

[^30]:    a W. H. Young (2).

[^31]:    a (2.8.4) of Theorem 13.

[^32]:    a It is modelled on the proof by F. Riesz (7) of the corresponding integral theorem: see § 6.8.

[^33]:    a Here, as in Theorem 26, we abandon our usual convention about $q$.

[^34]:    a We state the results for one variable and omit explicit reference to the range or set of integration.

[^35]:    a It is here that 'hompgeneity in $\Sigma$ ' (§ 1.4) is essential.
    b The precise form of argument used here was suggested to us by Mr H. D. Ursell.

[^36]:    a Or Buniakowsky's (see p. 16).

[^37]:    ${ }^{\text {a }}$ See § 5.2.

[^38]:    a $\int g^{k^{\prime}} d x=0$ would involve $g^{k^{\prime}} \equiv 0$ and so $g \equiv \infty$, and this possibility is excluded by the understanding of $\S 6.2$.
    ${ }^{\mathrm{b}}$ F. Riesz (2).

[^39]:    ${ }^{\text {a }}$ Continuity for $r<s$ can also be deduced from Theorems 111 and 197 (see § 6.12).
    b By § 6.3 (b) (i).
    c Strictly, the second case is excluded by the understanding of $\S$ 6.2.

[^40]:    a See Theorem 231.

[^41]:    a Extended to infinite series.
    ${ }^{b}$ A number of other analogues of theorems of Ch. III are stated among the miscellaneous theorems at the end of this chapter. A fuller treatment of some of them will be found in Jessen's papers 2 and 3. A good deal of the content of these papers has been incorporated, with the appropriate modifications, into Ch . III.

[^42]:    ${ }^{\text {a }}$ If, e.g., $b=\infty$, then $\beta=\lim _{x \rightarrow \infty} \phi(x)$.
    b If $g$ is any function of bounded variation, then $g=\phi-\psi$, where $\phi$ and $\psi$ are increasing functions, and we may define the Stieltjes integral of $f$ with respect to $g$ by $\int f d g=\int f d \phi-\int f d \psi$.

[^43]:    a We leave the specification of the exceptional cases to the reader.

[^44]:    a See for example de la Vallée Poussin (2), Hobson (1), Titchmarsh (1).
    ${ }^{b}$ This and the next example are due to Weierstrass and are of great historical importance.

[^45]:    a Euler's equation and Hilbert's invariant integral. Anything which we assume will be found without difficulty in the books of Bliss (1) or Bolza (1).

[^46]:    a By Theorem 181.

[^47]:    a Compare § 7.6.
    ${ }^{b}$ For which we are indebted to Prof. Bliss and Mr L. C. Young.

[^48]:    a The field differs slightly in character from those described in $\S 7.5$, since each extremal passes through $(0,0)$ and $(1,0)$.

[^49]:    a For the case $p=q=2$, see Hellinger and Toeplitz (1); for $q=p^{\prime}$, F. Riesz (1). The substance of the general theorems is to be found in M. Riesz (1). The important cases are naturally those in which $K$ has its least possible value, i.e. is the bound of $A$ (§8.8).
    ${ }^{\text {b }}$ Here $A \geqq 0$; but we write $|A|$ for $A$ in view of Theorem 288.

[^50]:    a Hellinger and Toeplitz (1), Schur (1).

[^51]:    a That is to say, assume non-negative values only for real $x$. The phrase 'positive form' has been used in this chapter in a different sense, that of a form with nonnegative coefficients and variables.

[^52]:    a All terms cancel except a number independent of $K$.
    ${ }^{b}$ The proof is that of Schur (1).

[^53]:    ${ }^{\text {a }}$ M. Riesz (1). The proof which we give of the theorem is substantially that of Riesz. An alternative proof has been given (not quite completely) by Paley $(2,4)$.

[^54]:    a Strictly, by Theorem 18 restated as Theorem 17 was restated in Theorem' 87.

[^55]:    a In this case $n=m$. A real unitary substitution is orthogonal.
    ${ }^{\text {b }}$ See §8.5. Riesz deduces these theorems in a different manner, and gives a number of other applications.

[^56]:    a That is to say it is true, for almost all $y$, that the two sides of (9.3.1) are equal for almost all $x$. See $\S 6.3(d)$.
    ${ }^{\mathrm{b}}$ Schur (1). Schur supposes $p=2$.

[^57]:    a For the case $p=2$ see Schur (1).
    ${ }^{\mathrm{b}}$ Cf. § 6.4; and see, for example, § 9.11.
    c Owen (1) gives a more general but less precise result.
    ${ }^{d}$ For the case $p=q=\ldots=r$ see Schur (1).

[^58]:    a. We describe Theorems 315 (together with the sharper Theorem 323) and $\mathbf{3 1 6}$ as 'Hilbert's theorems'. Strictly, Hilbert's theorem is Theorem 315, with $p=2$.
    b Associate elements of the integral symmetrically situated about the centre of the square of integration.

[^59]:    ${ }^{\text {a }}$ See F. Riesz (3), Hardy and Littlewood (2). We can express $f$ as the sum of two 'wurzelfrei' functions of $L$.

[^60]:    a Hardy and Littlewood (2). The theorem may also be stated in the form 'if a power series $g(z)=\Sigma b_{n} z^{n}$ is of bounded variation in $|z|<1$, then $\Sigma\left|b_{n}\right|$ is convergent'. For this form of the theorem, and for more precise results, see Fejér (1).
    b A much more general inequality, but without the best possible value of the constant, is proved by Hardy and Littlewood (1). See also Hardy (10).

[^61]:    ${ }^{\text {a }}$ It was a considerable time before any really simple proof of Hilbert's double series theorem was found.

[^62]:    a Partial integration is justified because $F^{p}$ is an integral when $F$ is an integral. For a formal proof see Hardy (4).

[^63]:    ${ }^{\text {a }}$ See § 10.17. Part of Theorem 329 is proved by Knopp (3).

[^64]:    a For a direct proof see Hardy (5).
    ${ }^{b}$ Copson (1); see also Hardy (6).
    c See § 8.7.

[^65]:    a For the details see Copson (1).
    b For the details (which are rather troublesome) see Hardy (4).
    c Compare $\S \S .4$ and 9.4 (1).

[^66]:    a Carleman (1). The proof given here is due to Pólya (2). The less precise convergence theorem (without the constant $e$ ) was found independently by other writers, and there are a number of proofs of one form or the other of the theorem. See Collingwood (in Valiron, 1, 186, where there is a proof due to Littlewood), Kaluza and Szegö (1), Knopp (1), Ostrowski (2, 201-204).

    $$
    \underset{n}{\Sigma}\left(a_{1} a_{2} \ldots a_{n}\right)^{1 / n} \leqq \sum_{n} \frac{1}{n} \sum_{m \leqq n} a_{m}=\Sigma a_{m} a_{n \leqq m} \sum_{n} \frac{1}{n} ;
    $$

    but the right-hand side is generally divergent. The proof fails because the $a$ in $a_{1} a_{2} \ldots a_{n}$ are 'too unequal', and too much is lost in replacing ${ }_{\infty}(a)$ by $\mathfrak{A}(a)$.

[^67]:    a Take $a_{1}=1, a_{2}=a_{3}=\ldots=0$. Then the result is false if $p>\frac{1}{2}$. For an alternative form of the result see Theorem 345.

[^68]:    a This theorem and Theorem 369 are valid for all real, not necessarily positive, $a$ and $b$.

[^69]:    a Theorem 43 (with $r=1$ and $p=1$ ) may be expressed, in our present notation, in the form

    $$
    n \Sigma \bar{a}_{j} \bar{b}_{n+1-j} \leqq \Sigma a_{j} \Sigma b_{j} \leqq n \Sigma \bar{a}_{j} \bar{b}_{j}
    $$

[^70]:    a $\mathfrak{a}^{i}$ means $\mathfrak{a}^{(i)}$ and $\alpha^{i}$ means $\alpha^{(i)}$ : in (10.3.1) above, $\mathfrak{a}^{2}$ is a power, but this use does not recur.
    ${ }^{\mathrm{b}} \mathfrak{a}^{1}$ is defined by this permutation.

[^71]:    a Hardy, Littlewood, and Pólya (1).
    b Hardy and Littlewood (4), Gabriel (1).
    c Gabriel (3).

[^72]:    a In order that the sets $\left(\mathfrak{c}^{\tau}\right)$ obtained by the process of $\S 10.3(2)$ should be symmetrical, in the sense of $\S 10.1$, we drop those $c^{\tau}$ which correspond to zero $\gamma^{\tau}$.
    ${ }^{\mathrm{b}}$ So that, strictly, we should write $\mathfrak{a}, \mathfrak{b}, \mathfrak{c}$ for $a, b, c$. There is, however, no further necessity for this notation.

[^73]:    a The argument is the same in principle as that used by Hardy, Littlewood, and Pólya (1), and substantially reproduced by Hardy and Littlewood (6). We have however expanded it considerably, Dr R. Rado having pointed out to us that the original form of the argument was not conclusive. Another form of the proof is indicated in Theorem 389 at the end of the chapter.

[^74]:    ${ }^{\text {a }}$ Gabriel (3). The case of the theorem in which all the sets are symmetrical was proved by Hardy and Littlewood (4).

[^75]:    a See Gabriel (3). A less precise inequality was given by Hardy and Littlewood (9).

[^76]:    a Compare the similar argument in § 10.6.
    b Unless the third function is nul.

[^77]:    a See for example de la Vallée Poussin (2, 20-23).
    b ${ }^{*} \chi(x)$ is naturally not to be confused with $\chi^{*}(x)$.

[^78]:    a See Fig. 3.
    ${ }^{b}$ We cannot lose two intervals because the intervals removed from $F_{x}$ and $G$ are the furthest to the right of their respective sets. This is the essential point of the proof.
    c The proof follows the line indicated by Zygmund (1). It is, however, considerably longer, and necessarily so, since Zygmund's proof is not, as it stands, conclusive.

    We proved that (10.15.16) is true when the intervals removed from $F$ and $G$ are

[^79]:    a See Fig. 4. In the figure the sets $E_{1}(t+h), \ldots$ are blackened on the sides of the triangle, the set $E_{123}(t+h)$ is shaded, the twelve perpendiculars are dotted, and the boundary of $E_{123}(t)-E_{123}(t+h)$ is indicated by a thick line.

[^80]:    a See Hardy and Littlewood (6).
    ${ }^{b}$ Due to F. Riesz (10).

[^81]:    a See the Miscellaneous Theorems at the end of the chapter.
    Theorem 385 was proved by Hardy and Littlewood (8), who deduced it by a limiting process from the analogous theorem for finite sums (Theorem 394). Their proof of Theorem 394 was elementary but long, and a much shorter proof was found by Gabriel (2). Riesz 'en combinant ce qui me paraît être l'idée essentielle de M. Gabriel avec un théorème appartenant aux éléments de l'analyse' (Lemma $A$ below) was able to prove the theorem directly and without limiting processes.

[^82]:    ${ }^{\text {a }}$ All that we need is that $g\left(\alpha_{k}\right) \leqq g\left(\beta_{k}\right)$; but the argument will probably be clearer if we make the lemma complete.

