## CHAPTER V

COMPOUND MATRICES
5.01 In chapter I it was found necessary to consider the adjoint of $A$ which is a matrix whose coordinates are the first minors of $|A|$. We shall now consider a more general class of matrices, called compound matrices, whose coordinates are minors of $|A|$ of the $r$ th order; before doing so, however, it is convenient to extend the definition of $S x y$ to apply to vectors of higher grade.
5.02 The scalar product Let $x_{i}=\Sigma \xi_{i j} e_{,}, y_{i}=\Sigma \eta_{i j} e_{j}(i=1,2, \cdots)$ be arbitrary vectors, then, by equation (37) §1.11 we have

$$
\begin{equation*}
\left|x_{1} x_{2} \cdots x_{r}\right|=\sum_{(i)}^{*}\left|\xi_{1 i_{1}} \xi_{2 i_{2}} \cdots \xi_{r i_{r}} \| e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right| \tag{1}
\end{equation*}
$$

and hence it is natural to extend the notion of the scalar product by setting
(2) $S\left|x_{1} x_{2} \cdots x_{r}\left\|y_{1} y_{2} \cdots y_{r}\left|=\sum_{(i)}^{*}\right| \xi_{1 i i} \xi_{2 i_{2}} \cdots \xi_{r i_{r}}\right\| \eta_{1 i_{1} \eta_{2 i_{2}}} \cdots \eta_{r i_{r}}\right|$.

We then have the following lemma which becomes the ordinary rule for multiplying together two determinants when $r=n$.

Lemma 1.

$$
\begin{equation*}
S\left|x_{1} x_{2} \cdots x_{z}\right|\left|y_{1} y_{2} \cdots y_{r}\right|=\left|S x_{i} y_{j}\right| \tag{3}
\end{equation*}
$$

For $S\left|x_{1} x_{2} \cdots x_{r}\right|\left|e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right|=\left|\xi_{1 i_{1} \xi_{2 i_{2}}} \cdots \xi_{r i_{r}}\right|$, hence

$$
\begin{aligned}
& S\left|x_{1} x_{2} \cdots x_{r}\right|\left|y_{1} e_{i_{2}} \cdots e_{i_{r}}\right|=\sum_{i_{1}} \eta_{1 i_{1}}\left|\xi_{1 i_{1}} \cdots \xi_{r i_{r}}\right| \\
& =\left|\left(\sum_{i_{1}} \eta_{1 i_{1}} \xi_{1 i_{1}}\right) \xi_{2 i_{2}} \cdots \xi_{r i_{r}}\right|=\left|S x_{1} y_{1}, \xi_{2 i_{2}} \cdots \xi_{r i_{r}}\right|
\end{aligned}
$$

again

$$
\begin{aligned}
S\left|x_{1} x_{2} \cdots x_{r}\right|\left|y_{1} y_{2} e_{i_{2}} \cdots e_{i_{r}}\right| & =\sum_{i_{2}} \eta_{2 i_{2}} S\left|x_{1} \cdots x_{r}\right|\left|y_{1} e_{i_{2}} \cdots e_{i_{r}}\right| \\
& =\sum_{i_{2}} \eta_{2 i_{2}}\left|S x_{1} y_{1} \xi_{2 i_{2}} \cdots \xi_{r i_{r}}\right| \\
& =\left|S x_{1} y_{1} S x_{2} y_{2} \xi_{3 i_{3}} \cdots \xi_{r i_{r}}\right|
\end{aligned}
$$

The lemma follows easily by a repetition of this process.
The Laplace expansion of a determinant can clearly be expressed as a scalar product. This is most easily done by introducing the notion of the comple-
ment of a vector relative to the fundamental basis. If $i_{1}, i_{2}, \cdots, i_{r}$ is a sequence of distinct integers in natural order each less than or equal to $n$ and $i_{r}+1, \cdots, i_{n}$ the remaining integers up to and including $n$, also arranged in natural order, the complement of $\left|e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right|$ relatively to the fundamental basis is defined as ${ }^{1}$

$$
\begin{equation*}
\left|e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right|_{c}=(-1)^{\Sigma i_{\alpha}+r(r+1) / 2}\left|e_{i_{r+1}} e_{i_{r+2}} \cdots e_{i_{n}}\right| \tag{4}
\end{equation*}
$$

and the complement of $\left|x_{1} x_{2} \cdots x_{r}\right|$ by

$$
\begin{equation*}
\left|x_{1} x_{2} \cdots x_{r}\right|_{c}=\sum_{(i)}^{*}\left|\xi_{1_{i}, \xi_{2 i_{2}}} \cdots \xi_{r i_{r}} \| e_{i_{1},} e_{i_{2}} \cdots e_{i_{r}}\right|_{c} \tag{5}
\end{equation*}
$$

which is a vector of grade $n-r$.
Laplace's expansion of a determinant in terms of minors of order $r$ can now be expressed in the following form.

Lemma 2.
(6) $S\left|x_{1} x_{2} \cdots x_{r}\right| c\left|x_{r+1} x_{r+2} \cdots x_{n}\right|=\left|\xi_{11} \xi_{22} \cdots \xi_{n n}\right|=\left|S x_{i} e_{j}\right|$
$=S\left|x_{1} \cdots x_{n}\left\|e_{1} \cdots e_{n}\left|=(-1)^{r(n-r)} S\right| x_{1} x_{2} \cdots x_{r}\right\| x_{r+1} \cdots x_{n}\right|_{c}$.
Further as an immediate consequence of (5) we have
Lemma 3.

$$
\begin{equation*}
S\left|x_{1} x_{2} \cdots x_{r}\right|_{c}\left|y_{1} y_{2} \cdots y_{r}\right|_{c}=S\left|x_{1} x_{2} \cdots x_{r}\right|\left|y_{1} y_{2} \cdots y_{r}\right| . \tag{7}
\end{equation*}
$$

5.03 Compound matrices. If $A=\Sigma a_{i,}, e_{i j}$, then, as in (1),

$$
\left|A x_{1} A x_{2} \cdots A x_{r}\right|=\sum_{(j)}^{*}\left|\xi_{1_{1}} \cdots \xi_{r_{r}} \| A e_{\mu_{1}} \cdots A e_{i_{r}}\right| .
$$

But $A e_{j}=\sum_{i} a_{i j} e_{i}$; so a second application of (1) gives

$$
\left|A x_{1} A x_{2} \cdots A x_{r}\right|=\sum_{(i)}^{*} \sum_{(j)}^{*}\left|\xi_{11_{1}} \cdots \xi_{r j_{r}}\left\|a_{i, i_{1}} \cdots a_{i r_{r} r}\right\| e_{i_{1}} \cdots e_{i_{r}}\right|
$$

But the determinants $\left|\xi_{11_{1}} \cdots \xi_{r j_{r}}\right|$ are the coordinates of the $r$-vector $\left|x_{1} x_{2} \cdots x_{r}\right|$; hence $\left|A x_{1} \cdots A x_{r}\right|$ is a linear vector form in $\left|x_{1} x_{2} \cdots x_{r}\right|$ in the corresponding space of $\binom{n}{r}$ dimensions. We denote this vector function or matrix by $C_{r}(A)$ and write

$$
\begin{equation*}
\left|A x_{1} A x_{2} \cdots A x_{r}\right|=C_{r}(A)\left|x_{1} x_{2} \cdots x_{r}\right| . \tag{8}
\end{equation*}
$$

We shall call $C_{r}(A)$ the $r$ th compound of $A$. Important particular cases are

$$
C_{3}(A)=A, \quad C_{n}(A)=|A|
$$

${ }^{1}$ The Grassmann notation cannot be conveniently used here since it conflicts with the notation for a determinant. It is sometimes convenient to define the complement of $\left|e_{1} e_{2} \cdots e_{n}\right|$ as 1 .
and, if $k$ is a scalar, ( $8^{\prime \prime}$ )

$$
C_{r}(k)=k^{r} .
$$

Theorem 1.

$$
\begin{equation*}
C_{r}(A B)=C_{r}(A) C_{r}(B) \tag{9}
\end{equation*}
$$

For

$$
\begin{aligned}
\left|A B x_{1} A B x_{2} \cdots A B x_{r}\right| & =C_{r}(A)\left|B x_{1} B x_{2} \cdots B x_{r}\right| \\
& =C_{r}(A) C_{r}(B)\left|x_{1} x_{2} \cdots x_{r}\right| .
\end{aligned}
$$

Corollary. If $|A| \neq 0$, then

$$
\begin{equation*}
\left[C_{r}(A)\right]^{-1}=C_{r}\left(A^{-1}\right) . \tag{10}
\end{equation*}
$$

Theorem 2.

$$
\begin{equation*}
\left[C_{r}(A)\right]^{\prime}=C_{r}\left(A^{\prime}\right) . \tag{11}
\end{equation*}
$$

For $S\left|x_{1} x_{2} \cdots x_{r}\right| C_{r}(A)\left|y_{1} y_{2} \cdots y_{r}\right|=\left|S x_{i} A y_{j}\right|=\left|S A^{\prime} x_{i} y_{j}\right|$

$$
=S\left|A^{\prime} x_{1} \cdots A^{\prime} x_{r} \| y_{1} \cdots y_{r}\right|=S\left|y_{1} \cdots y_{r}\right| C_{r}\left(A^{\prime}\right)\left|x_{1} \cdots x_{r}\right|
$$

Theorem 3. If $A=\sum_{i}^{m} a_{i} S b_{i}$, then

$$
\begin{equation*}
C_{r}(A)=\sum_{(i)}^{*}\left|a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}\right| S\left|b_{i_{1}} b_{i_{2}} \cdots b_{i_{r}}\right| \tag{12}
\end{equation*}
$$

This theorem follows by direct substitution for $A$ ir the left-hand side of (8) It gives a second proof for Theorem 2.

If $r=m$, (12) consists of one term only, and this term is 0 unless $m$ is the rank of $A$, a property which might have been made the basis of the definition of rank. In particular, if $X=\sum_{1}^{r} e_{i} S x_{i}, Y=\sum_{1}^{r} y_{i} S e_{i}$, then $C_{r}(X)=$ $\left|e_{1} e_{2} \cdots e_{r}\right| S\left|x_{1} x_{2} \cdots x_{r}\right|, C_{r}(Y)=\left|y_{1} y_{2} \cdots y_{r}\right| S\left|e_{1} e_{2} \cdots e_{r}\right|$ so that $C_{r}(X Y)=\left|e_{1} e_{2} \cdots e_{r}\right| S\left|x_{1} x_{2} \cdots x_{r}\right|\left|y_{1} y_{2} \cdots y_{r}\right| S\left|e_{1} e_{2} \cdots e_{r}\right|$. But $X Y=\sum_{i, j} e_{i} S x_{i} y_{j} S e_{i}$ so that $C_{r}(X Y)=\left|S x_{i} y_{i} \| e_{1} c_{2} \cdots e_{r}\right| S\left|e_{1} e_{2} \cdots e_{r}\right|$. Comparing these two forms of $C_{r}(X Y)$ therefore gives another proof of the first lemma of §5.02.

If we consider the complement of $\left|A x_{1} A x_{2} \cdots A x_{r}\right|$ we arrive at a new matrix $C^{r}(A)$ of order $\binom{n}{r}$ which is called the $r$ th supplementary compound of $A$. From (7) and (12) we have

$$
\begin{align*}
\left|A x_{1} A x_{2} \cdots A x_{r}\right|_{c} & =\sum_{i}^{*}\left|a_{i_{1}} \cdots a_{i_{r}}\right|_{c} S\left|b_{i_{1}} \cdots b_{i_{r}}\right|_{c}\left|x_{1} \cdots x_{r}\right|_{c}  \tag{13}\\
& =C^{r}(A)\left|x_{1} x_{2} \cdots x_{r}\right|_{c}
\end{align*}
$$

and derive immediately the following which are analogous to Theorems 1 and 2.

Theorem 4.

$$
\begin{equation*}
C^{r}(A B)=C^{r}(A) C^{r}(B) \tag{14}
\end{equation*}
$$

Theorem 5.

$$
\begin{equation*}
\left[C^{r}(A)\right]^{\prime}=C^{r}\left(A^{\prime}\right) \tag{15}
\end{equation*}
$$

The following theorems give the connection between compounds and supplementary compounds and also compounds of compounds.

Theorem 6.

$$
\begin{equation*}
C^{r}\left(A^{\prime}\right) C_{n-r}(A)=|A|=C^{n-r}(A) C_{r}\left(A^{\prime}\right) \tag{16}
\end{equation*}
$$

This is the Laplace expansion of the determinant $|A|$. Using equation (6) and setting $|e|$ for $\left|e_{1} e_{2} \cdots e_{n}\right|$ we have

$$
\begin{aligned}
|A| S \mid & \left.x_{1} x_{2} \cdots x_{r}\right|_{c}\left|x_{r+1} \cdots x_{n}\right|=|A| S\left|x_{1} \cdots x_{n}\right||e| \\
& =S\left|A x_{1} \cdots A x_{n}\right||e| \\
& =S\left|A x_{1} \cdots A x_{r}\right|_{c}\left|A x_{r+1} \cdots A x_{n}\right| \\
& =S C^{r}(A)\left|x_{1} \cdots x_{r}\right|_{c} C_{n-r}(A)\left|x_{r+1} \cdots x_{n}\right| \\
& =S\left|x_{1} \cdots x_{r}\right|_{c} C^{r}\left(A^{\prime}\right) C_{n-r}(A)\left|x_{r+1} \cdots x_{n}\right|
\end{aligned}
$$

and, since the $x$ 's are arbitrary, the first part of the theorem follows. The second part is proved in a similar fashion.

Putting $r=n-1$ in (16) gives the following corollary.
Corollary.

$$
\operatorname{adj} A=C^{n-1}\left(A^{\prime}\right)
$$

Theorem 7.

$$
\begin{equation*}
\left|C_{r}(A)\right|=|A|^{\binom{n-1}{r-1}}=\left|C^{r}(A)\right| . \tag{17}
\end{equation*}
$$

For from (16) with $A^{\prime}$ in place of $A$, and from the fact that the order of $C_{r}(A)$ is $\binom{n}{r}$, we have

$$
|A|^{\binom{n}{r}}=\left|C^{r}(A) C_{n-r}\left(A^{\prime}\right)\right|=\left|C^{r}(A)\right|\left|C_{n-r}\left(A^{\prime}\right)\right|
$$

and, since $|A|$ is irreducible when the coordinates of $A$ are arbitrary variables, it follows that $\left|C^{r}(A)\right|$ is a power of $|A|$. Considerations of degree then show that the theorem is true when the coordinates are variables and, since the identity is integral, it follows that it is also true for any particular values of these variables.

Theorem 8.

$$
\begin{align*}
& |A|^{\binom{n-1}{r}} C_{s}\left(C_{r}(A)\right)=|A|^{s} C^{\binom{n}{r}-s}\left(C^{n-r}(A)\right)  \tag{18}\\
& |A|^{\binom{n-1}{r}} C_{s}\left(C^{r}(A)\right)=|A|^{s} C^{\binom{n}{r}-s}\left(C_{n-r}(A)\right) . \tag{19}
\end{align*}
$$

Using (15), (16) and (17) we get

$$
C_{s}\left(C^{n-r}\left(A^{\prime}\right)\right) C^{\binom{n}{r}-s}\left(C^{n-r}(A)\right)=\left|C^{n-r}(A)\right|=|A|^{\binom{n-1}{r}}
$$

therefore

$$
\begin{aligned}
|A|^{\left(n_{r}^{1}\right)} C_{s}\left(C_{r}(A)\right) & =C_{s}\left(C_{r}(A)\right) C_{s}\left(C^{n-r}\left(A^{\prime}\right)\right) C^{n-s}\left(C^{n-r}(A)\right) \\
& =C_{s}\left(C_{r}(A) C^{n-r}\left(A^{\prime}\right)\right) C^{n-s}\left(C^{n-r}(A)\right) \\
& =C_{s}(|A|) C^{n-s}\left(C^{n-r}(A)\right) \\
& =|A|^{\bullet} C^{n-s}\left(C^{n-r}(A)\right) .
\end{aligned}
$$

Similarly

$$
C_{s}\left(C_{n-r}\left(A^{\prime}\right)\right) C_{s}\left(C^{r}(A)\right)=C_{s}(|A|)=|A|^{s}
$$

and therefore

$$
\begin{aligned}
|A|^{s} C^{\binom{n}{r}-s}\left(C_{n-r}(A)\right) & =C^{\binom{n}{r}-s}\left(C_{n-r}(A)\right) C_{s}\left(C_{n-r}\left(A^{\prime}\right)\right) C_{s}\left(C^{r}(A)\right) \\
& =\left|C_{n-r}(A)\right| C_{s}\left(C^{r}(A)\right) \\
& =|A|^{\binom{n-1}{r}} C_{s}\left(C^{\prime}(A)\right) .
\end{aligned}
$$

An important particular case is $C_{s}\left(C^{n-1}(A)\right)=|A|^{\theta^{-1} C^{n-}}(A)$ whence

$$
\begin{equation*}
C_{s}(\operatorname{adj} A)=C_{s}\left(C^{n-1}\left(A^{\prime}\right)\right)=|A|^{s-1} C^{n-s}\left(A^{\prime}\right) . \tag{20}
\end{equation*}
$$

5.04 Roots of compound matrices. If $A$ has simple elementary divisors and its roots are $\lambda_{1}, \lambda_{2}, \cdots, \lambda_{n}$, the corresponding invariant vectors being $a_{1} . a_{2}, \cdots, a_{n}$, then the roots of $C_{r}(A)$ are the products $\lambda_{i_{1} \lambda_{i_{2}}} \cdots \lambda_{i_{r}}$ in which no two subscripts are the same and the subscripts are arranged in, say, numerical order; and the invariant vector corresponding to $\lambda_{i_{1}} \lambda_{i_{2}} \cdots \lambda_{i_{r}}$ is $\left|a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}\right|$. For there are $\binom{n}{r}$ distinct vectors of this type and

$$
C_{r}(A)\left|a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}\right|=\left|A a_{i_{1}} \cdots A a_{i_{r}}\right|=\lambda_{i_{1}} \lambda_{i_{2}} \cdots \lambda_{i_{r}}\left|a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}\right|
$$

Similarly for $C^{r}(A)$ the roots and invariants are $\lambda_{i_{1}} \lambda_{i_{2}} \cdots \lambda_{i_{r}}$ and $\left|a_{i_{1}} a_{i_{2}} \cdots a_{i_{r}}\right|_{c}$.

It follows from considerations of continuity that the roots are as given above even when the elementary divisors are not simple.
5.05 Bordered determinants. Let $A=\left\|a_{i j}\right\|=\sum_{j=1}^{n} a_{j} S e_{1}, a_{i}=\sum_{i} a_{i j} e_{i}$, be any matrix and associate with it two sets of vectors

$$
\begin{aligned}
& X: x_{i}=\sum_{j=1}^{n} \xi_{i j} e_{j}, \\
& Y: y_{i}=\sum_{j=1}^{n} \eta_{i j} e_{j} .
\end{aligned}
$$

Consider the bordered determinant
where $r<n$, and $0_{r}$ is a square block of 0 's with $r$ rows and columns.
If we introduce $r$ additional fundamental units $e_{n+1}, \cdots, e_{n+r}, \Delta_{r}$ can be regarded as the determinant of a matrix $\mathfrak{A}$ of order $n+r$, namely,

$$
\mathfrak{A}=\sum_{1}^{n} a_{i} S e_{i}+\sum_{1}^{r} x_{i} S e_{n+i}+\sum_{i}^{r} e_{n+i} S y_{i}=\sum_{1}^{n+2 r} c_{i} S d_{i} .
$$

If now we form $|\mathfrak{A}|=S|e| C_{n+r}(\mathfrak{A})|e|$ as in $\S 5.03$, we have

$$
C_{n+r}(\mathfrak{H})=\sum_{(i)}^{*}\left|c_{i_{1}} \cdots c_{i_{n+r}}\right| S\left|d_{i_{1}} \cdots d_{i_{n-r}}\right|(i=1,2, \cdots, n+2 r)
$$

In this form any $\left|c_{i_{1}} \cdots c_{i_{n+r}}\right|$ which contains more than $n$ out of $a_{1}, \cdots, a_{n}$. $x_{1}, \cdots, x_{r}$ is necessarily 0 ; also, if it does not contain all the $x$ 's, the corresponding $\left|d_{i_{1}}, \cdots, d_{i_{n+r}}\right|$ will contain more than $n$ out of $e_{1}, \cdots, e_{n}$, $y_{1}, \cdots, y_{r}$ and is consequently 0 . We therefore have

$$
\begin{aligned}
& C_{n+r}(\mathfrak{V})=\sum_{(i)}^{*}\left|a_{i_{1}} a_{i_{2}} \cdots a_{i_{n-r}} x_{1} x_{2} \cdots x_{r} e_{n+1} \cdots e_{n+r}\right| \\
& \times S\left|e_{i_{1} e_{i_{2}}} \cdots e_{i_{n-r}} y_{1} y_{2} \cdots y_{r} e_{n+1} \cdots e_{n+r}\right| \quad(i=1,2, \cdots, n)
\end{aligned}
$$

and hence, passing back to space of $n$ dimensions,

$$
\begin{aligned}
|\mathfrak{A}| & =\sum_{i}^{*} S\left|e\left\|a_{i_{1}} \cdots a_{i_{n-r}} x_{1} \cdots x_{r}|S| e_{i_{1}} \cdots e_{i_{n-r}} y_{1} \cdots y_{r}\right\| e\right| \\
& =\Sigma^{*} S\left|x_{1} \cdots x_{r} \| a_{i_{1}} \cdots a_{i_{n-r}}\right| c S\left|e_{i_{1}} \cdots e_{i_{n-r}}\right| c\left|y_{1} \cdots y_{r}\right| \\
& =S\left|x_{1} \cdots x_{r}\right| C^{n-r}(A)\left|y_{1} \cdots y_{r}\right| .
\end{aligned}
$$

This relation shows why the bordered determinant is frequently used in place of the corresponding compound in dealing with the theory of forms.
5.06 The reduction of bilinear forms. The Lagrange method of reducing quadratic and bilinear forms to a normal form is, as we shall now see, closely connected with compounds.

If $A$ is any matrix, not identically 0 , there exist vectors $x_{1}, y_{1}$ such that $S x_{1} A y_{1} \neq 0$; then, setting $A=A_{1}$ for convenience, the matrix

$$
A_{2}=A_{1}-A_{1} y_{1} \frac{S A_{1}^{\prime} x_{1}}{S x_{1} A_{1} y_{1}}
$$

has its rank exactly 1 less than that of $A$. For, if $A_{1} z=0$, then

$$
A_{2} z=A_{1} z-A_{1} y_{1} \frac{S A^{\prime} x_{1} \cdot z}{S x_{1} A_{1} y_{1}}=A_{1} z-A_{1} y_{1} \frac{S x_{1} A_{1} z}{S x_{1} A_{1} y_{1}}=0
$$

and, conversely if $A_{2} z=0$, then

$$
A_{1} z=A_{1} y_{1} \frac{S x_{1} A_{1} z}{S x_{1} A_{1} y_{1}}=k A_{1} y_{1}
$$

say, or $A_{1}\left(z-k y_{1}\right)=0$. The null-space of $A_{2}$ is therefore obtained from that of $A_{1}$ by adding $y_{1}$ to its basis, which increases the order of this space by 1 since $A_{1} y_{1} \neq 0$.

If $A_{2} \neq 0$, this process may be repeated, that is, there exist $x_{2}, y_{2}$ such that $S x_{2} A_{2} y_{2} \neq 0$ and the rank of

$$
A_{3}=A_{2}-A_{2} y_{2} \frac{S A_{2}^{\prime} x_{2}}{S x_{2} A_{2} y_{2}}
$$

is 1 less than that of $A_{2}$. If the rank of $A$ is $r$, we may continue this process by setting

$$
\begin{equation*}
A_{s+1}=A_{s}-A_{s} y_{s} \frac{S A_{s}^{\prime} x_{s}}{S x_{s} A_{s} y_{s}} \quad(s=1,2, \cdots, r) \tag{22}
\end{equation*}
$$

where $S x_{s} A_{s} y_{s} \neq 0$ and $A_{1}=A, A_{r+1}=0$; we then have

$$
\begin{equation*}
A=\sum_{s=1}^{r} A_{s} y_{s} \frac{S A_{s}{ }_{s} x_{s}}{S x_{s} A_{s} y_{s}}=\sum_{1}^{r} \mathfrak{A}_{s} \tag{23}
\end{equation*}
$$

where $\mathfrak{H}_{s}=A_{s} y_{s} \frac{S A_{s}^{\prime} x_{s}}{S x_{s} A_{s} y_{s}}$ is a matrix of rank 1. Generally speaking, one may take $x_{s}=y_{s}$ and it is of some interest to determine when this is not possible. If $S x B x=0$ for every $x$, we readily see that $B$ is skew. For then $S e_{i} B e_{i}$ $=S e_{j} B e_{j}=S\left(e_{i}+e_{j}\right) B\left(e_{i}+e_{2}\right)=0$ and therefore

$$
0=S\left(e_{i}+e_{\jmath}\right) B\left(e_{i}+e_{\jmath}\right)=S e_{i} B e_{i}+S e_{j} B e_{j}+S e_{i} B e_{\jmath}+S e_{\jmath} B e_{i},
$$

that is, $S e_{i} B e_{\jmath}=-S e_{j} B e_{i}$ and hence $B^{\prime}=-B$. Hence we may take $x_{s}=y_{s}$ so long as $A_{s} \neq-A_{s}^{\prime}$.
5.07 We shall now derive more explicit forms for the terms in (23) and show how they lead to the Sylvester-Francke theorems on compound determinants.

Let $x^{1}, x^{2}, \cdots, x^{r}, y^{1}, y^{2}, \cdots, y^{r}$ be variable vectors and set

$$
\begin{align*}
J & =S\left|x_{s} x^{1} x^{2} \cdots x^{r}\right| C_{r+1}\left(A_{s}\right)\left|y_{s} y^{1} y^{2} \cdots y^{r}\right|  \tag{24}\\
& =S\left|x_{s} x^{1} x^{2} \cdots x^{r}\right|\left|A_{s} y_{s} A_{s} y^{1} \cdots A_{s} y^{r}\right| ;
\end{align*}
$$

then from (22)

$$
\begin{aligned}
J & =S\left|x_{s} x^{1} \cdots x^{r}\right|\left|A_{s} y_{s} A_{s+1} y^{1} \cdots A_{s}+{ }_{1} y^{r}\right| \\
& =\left|S x_{s} A_{s} y_{s} S x^{1} A_{s+1} y^{1} \cdots S x^{r} A_{s+1} y^{r}\right| .
\end{aligned}
$$

If the $x$ 's denote rows in this determinant, the first row is

$$
S x_{s} A_{s} y_{s}, S x_{s} A_{s+1} y^{1}, \cdots, S x_{s} A_{s+1} y^{r}
$$

each term of which is 0 except the first, since $x_{s}$ lies in the null-space of $A_{s+1}^{\prime}$, and $S x_{s} A_{s} y_{s} \neq 0$. Hence

$$
\begin{equation*}
J=S x_{s} A_{s} y_{s}\left|S x^{1} A_{s+1} y^{1} \cdots S x^{r} A_{s+1} y^{r}\right| \tag{25}
\end{equation*}
$$

and therefore from (24)

$$
\begin{gather*}
S\left|x_{s} x^{1} \cdots x^{r}\right| C_{r+1}\left(A_{s}\right)\left|y_{s} y^{1} \cdots y^{r}\right|  \tag{26}\\
=S x_{s} A_{s} y_{s} S\left|x^{1} \cdots x^{r}\right| C_{r}\left(A_{s+1}\right)\left|y^{1} \cdots y^{r}\right| .
\end{gather*}
$$

Repeated application of this relation gives

$$
\begin{gather*}
S\left|x_{s} x_{s+1} \cdots x_{s+t-1} x^{1} x^{2} \cdots x^{r}\right| C_{r+t}\left(A_{s}\right)\left|y_{s} \cdots y_{s+t-1} y^{1} \cdots y^{r}\right|  \tag{27}\\
=S x_{s} A_{s} y_{s} S x_{s+1} A_{s+1} y_{s+1} \cdots S x_{s+t-1} A_{s+t-1} y_{s+t-1} S\left|x^{1} \cdots x^{r}\right| \\
\cdot C_{r}\left(A_{s+t}\right)\left|y^{1} \cdots y^{r}\right|,
\end{gather*}
$$

a particular case of which is

$$
\begin{align*}
S \mid x_{1} x_{2} \cdots & x_{s-1} x\left|C_{s}(A)\right| y_{1} \cdots y_{s-1} y \mid  \tag{27'}\\
& =S x_{1} A_{1} y_{1} \cdots S x_{s-1} A_{s-1} y_{s-1} S x A_{s} y .
\end{align*}
$$

To simplify these and similar formulae we shall now use a single letter to indicate a sequence of vectors; thus we shall set $X_{s, s+t-1}$ for $x_{s} x_{s+1} \ldots$ $x_{s+t-1}$ and $Y^{r}$ for $y^{1} y^{2} \cdots y^{r}$; also $C_{r, ~ s}$ for $C_{r}\left(A_{s}\right)$. Equations (26) and (27) may then be written

$$
\begin{gather*}
S\left|x_{s} X^{r}\right| C_{r+1, s}\left|y_{s} Y^{r}\right|=S x_{s} A_{s} y_{s} S\left|X^{r}\right| C_{r, s+1}\left|Y^{r}\right|,  \tag{26a}\\
S\left|X_{s, s+t-1} X^{r}\right| C_{r+t, s}\left|Y_{s, s+t-1} Y^{r}\right|=\prod_{i=s}^{0+-1} S x_{i} A_{i} y_{i} S\left|X^{r}\right| C_{r, s+t}\left|Y^{r}\right| \tag{27a}
\end{gather*}
$$

We get a more convenient form for (26a), namely

$$
\begin{gather*}
S\left|X_{s, t} X^{r}\right| C_{r+t-s+1, s}\left|Y_{s, t} Y^{r}\right|  \tag{28}\\
=S x_{s} A_{s} y_{s} S\left|X_{s+1, t} X^{r}\right| C_{r+t-s, s+1}\left|Y_{s+1, t} Y^{r}\right|
\end{gather*}
$$

by replacing $r$ by $r+t-s$ and then changing $x^{1} x^{2} \cdots x^{r+t-s}$ into $x_{s+1} \cdots x_{t} x^{1} \cdots x^{r}$ along with a similar change in the $y$ 's. Putting $s$ $=1,2, \cdots, i$ in succession and forming the product of corresponding sides of
the equations so obtained from (28) we get after canceling the common factors, which are not identically 0 provided that $r+t$ is not greater than the rank of $A$,

$$
\begin{equation*}
S\left|X_{t} X^{r}\right| U_{r+t, 1}^{\prime}\left|Y_{t} Y^{r}\right|=\prod_{1}^{t} S x_{i} A_{i} y_{i} \cdot S\left|X^{r}\right| C_{r, t+1}\left|Y^{r}\right| \tag{29}
\end{equation*}
$$

or rrom (27')

$$
\begin{equation*}
S\left|X_{t} X^{r}\right| C_{r+t}\left|Y_{t} Y^{r}\right|=S\left|X_{t}\right| C_{t}\left|Y_{t}\right| S\left|X^{r}\right| C_{r, t+1}\left|Y^{r}\right| \tag{30}
\end{equation*}
$$

which may also be written in the form

$$
K \equiv \frac{S\left|X_{t} X^{r}\right| C_{r+s}\left|Y_{t} Y^{r}\right|}{S\left|X_{t}\right| C_{t}\left|Y_{t}\right|}=\left|S x^{i} A_{t+1} y^{i}\right|
$$

in particular

$$
\begin{equation*}
\frac{S\left|X_{t} x\right| C_{t+1}(A)\left|Y_{t} y\right|}{S\left|X_{t}\right| C_{t}(A)\left|Y_{t}\right|}=S x A_{t+1} y \tag{31}
\end{equation*}
$$

This gives a definition of $A_{t+1}$ which may be used in place of (22); it shows that this matrix depends on $2 t$ vector parameters. It is more convenient for some purposes to use the matrix $A^{(t)}$ defined by

$$
\begin{equation*}
S x A^{(t)} y=S\left|X_{t} \dot{x}\right| C_{t+1}(A)\left|Y_{t} y\right| \tag{32}
\end{equation*}
$$

From (31) we then have $S x^{i} A_{t+1} y^{j}=S x^{i} A^{(t)} y^{j} / S\left|X_{t}\right| C_{t}\left|Y_{t}\right|$ and therefore from ( $30^{\prime}$ )

$$
\begin{equation*}
K=\frac{\left|S x^{i} A^{(t)} y^{j}\right|}{\left[S\left|X_{t}\right| C_{t}\left|Y_{t}\right|\right]^{r}}=\frac{S\left|X^{r}\right| C_{r}\left(A^{(t)}\right)\left|Y^{r}\right|}{\left[S\left|X_{t}\right| C_{t}(A)\left|Y_{t}\right|\right]^{r}} \tag{33}
\end{equation*}
$$

Hence

$$
\begin{equation*}
S\left|X_{t} X^{r}\right| C_{r+t}(A)\left|Y_{t} Y^{r}\right|=\frac{S\left|X^{r}\right| C_{r}\left(A^{(t)}\right)\left|Y^{r}\right|}{\left[S\left|X_{t}\right| C_{t}(A)\left|Y_{t}\right|\right]^{r-1}} \tag{34}
\end{equation*}
$$

which is readily recognized as Sylvester's theorem if the $x$ 's are replaced by fundamental units and the integral form of (33) is used.
5.08 Invariant factors. We shall now apply the above results in deriving the normal form of $\S 3.02$. We require first, however, the following lemma.

Lemma 4. If $A(\lambda)$ is a matric polynomial, there exists a constant vector $y$ and a vector polynomial $x$ such that $S x A y$ is the highest common factor of the coordinates of $A$.

Let $y=\Sigma \eta_{i} e_{i}$ be a vector whose coordinates are variables independent of $\lambda$.
Let $\alpha_{1}$ be the H. C. F. of the coordinates of $A=\left\|a_{i},\right\|$ and set

$$
A=\alpha_{1} B, B y=\Sigma \eta, b_{i}, e_{1}=\Sigma \beta_{i} e_{i}
$$

There is no value $\lambda_{1}$ of $\lambda$ independent of the $\eta$ 's for which every $\beta_{i}=0$; for if this were so, $\lambda-\lambda_{1}$ would be a factor of each $b_{i}$, and $\alpha_{1}$ could not then be the H. C. F. of the $a_{i r}$. Hence the resultant of $\beta_{1}, \beta_{2}, \cdots, \beta_{n}$ as polynomials in $\lambda$ is not identically 0 as a polynomial in the $\eta$ 's; there are therefore values of the $\dot{\eta}$ 's for which this resultant is not 0 , and for these values the $\beta$ 's have no factor common to all. There then exist scalar polynomials $\xi_{1}, \xi_{2}, \cdots, \xi_{n}$ such that $\Sigma \xi_{i} \beta_{i}=1$ and therefore, if $x=\Sigma \xi_{i} e_{i}$, we have $S x B y=1$ or $S x A y=\alpha_{1}$.

Returning now to the form of $A$ given in $\S 5.06$, namely

$$
A=\sum_{1}^{r} \frac{A_{s} y_{s} S A_{s}^{\prime} x_{s}}{S x_{s} A_{s} y_{s}},
$$

we can as above choose $x_{s}, y_{s}$ in such a manner that $S x_{s} A_{s} y_{s}=\alpha_{s}$ is the highest common factor of the coordinates of $A_{s}$ and, when this is done, $v_{s}=A_{s} y_{s} / \alpha_{s}$, $u_{s}=A_{s}^{\prime} x_{s} / \alpha_{s}$ are integral in $\lambda$. We then have

$$
\begin{equation*}
A=\sum_{1}^{r} \frac{A_{s} y_{s} S A_{s}^{\prime} x_{s}}{\alpha_{s}}=\Sigma \alpha_{s} v_{s} S u_{s} \tag{35}
\end{equation*}
$$

Moreover $A_{s} y_{i}=\mathbf{0}=A_{s}^{\prime} x_{i}$ when $i<s$ and therefore in

$$
S\left|x_{1} \cdots x_{r}\right|\left|A_{1} y_{1} A_{2} y_{2} \cdots A_{r} y_{r}\right|=\left|S x_{i} A_{j} y_{j}\right|=\left|S A_{j}^{\prime} x_{i} y_{2}\right|
$$

all terms on one side of the main diagonal are 0 so that it reduces to $S x_{1} A_{1} y_{1}$ $\cdots S x_{r} A_{r} y_{r}=\alpha_{1} \alpha_{2} \cdots \alpha_{r}$. Hence, dividing by $\alpha_{1} \cdots \alpha_{r}$ and replacing $A_{j} y_{i} / \alpha_{j}$ by $v_{j}$ as above, we see that $\left|x_{1} \cdots x_{r}\right|$ and $\left|v_{1} \cdots v_{r}\right|$ are not 0 for any value of $\lambda$, and therefore the constituent vectors in each set remain linearly independent for all values of $\lambda$. It follows in the same way that the sets $u_{1}, \cdots, u_{r}$ and $y_{1}, \cdots, y_{r}$, respectively, are also linearly independent for all values of $\lambda$, that is, these four sets are elementary sets. It follows from Theorem 5 §4.03, that we can find elementary polynomials $P$ and $Q$ such that

$$
P v_{i}=e_{i}=Q^{\prime} u_{i} \quad(i=1,2, \cdots, r)
$$

and hence

$$
\begin{equation*}
P A Q=P\left(\sum_{1}^{r} \alpha_{s} v_{s} S u_{s}\right) Q=\sum_{1}^{r} \alpha_{s} e_{s} S e_{s}, \tag{36}
\end{equation*}
$$

which is the normal form of $\S 3.02$.
5.09 Vector products. Let $x_{i}=\Sigma \xi_{i j} e,(i=1,2, \cdots, r)$ be a set of arbitrary vectors and consider the set of all products of the form $\xi_{1 i_{1}} \xi_{2_{i}} \cdots \xi_{r i_{r}}$ arranged in some definite order. These products may then be regarded as the coordinates of a hypernumber ${ }^{2}$ of order $n^{r}$ which we shall call the tensor product ${ }^{3}$

[^0]of $x_{1}, x_{2}, \cdots, x_{r}$ and we shall denote it by $x_{1} x_{2} \cdots x_{r}$. In particular if we take all the products $e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\left(i_{1}, i_{2}, \cdots, i_{r}=1,2, \cdots, n\right)$ each has all its coordinates zero except one, which has the value 1 , and no two are equal. Further
$$
x_{1} x_{2} \cdots x_{r}=\Sigma \xi_{11_{1},} \xi_{2 i_{2}} \cdots \xi_{r i_{r}} e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}
$$

If we regard the products $e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}$ as the basis of the set of hypernumbers, we are naturally led to consider sums of the type .

$$
w=\Sigma \omega_{i_{1} i_{2}} \cdots i_{r} e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}
$$

where the $\omega$ 's are scalars; and we shall call such a hypernumber a tensor of grade $r$. It is readily seen that the product $x_{1} x_{2} \cdots x_{r}$ is distributive and homogeneous with regard to each of its factors, that is,

$$
x_{1}\left(\lambda x_{2}+\mu y_{2}\right) x_{3} \cdots x_{r}=\lambda x_{1} x_{2} \cdots x_{r}+\mu x_{1} y_{2} x_{3} \cdots x_{r} .
$$

The product of two tensors of grade $r$ and $s$ is then defined by assuming the distributive law and setting

$$
\left(e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right)\left(e_{\lambda_{1}} e_{j_{2}} \cdots e_{j_{s}}\right)=e_{i_{1}} \cdots e_{i_{r}} e_{2_{1}} \cdots e_{j_{s}}
$$

It is easily shown that the product so defined is associative; it is however not commutative as is seen from the example

$$
\begin{aligned}
x_{1} x_{2}-x_{2} x_{1} & =\Sigma \Sigma\left(\xi_{1 i_{1}} \xi_{2 i_{2}}-\xi_{1 i_{2}} \xi_{2 i_{1}}\right) e_{i_{1}} e_{i_{2}} \\
& =\sum_{(i)}^{*}\left|\begin{array}{cc}
\xi_{1 i_{1}} & \xi_{1 i_{2}} \\
\xi_{2 i_{1}} & \xi_{2 i_{2}}
\end{array}\right|\left(e_{i_{1}} e_{i_{2}}-e_{i_{2}} e_{i_{1}}\right) .
\end{aligned}
$$

Here the coefficients of $e_{i_{1}} e_{i_{2}}-e_{i_{2}} e_{i_{1}}\left(i_{1}<i_{2}\right)$ are the coordinates of $\left|x_{1} x_{2}\right|$ so that this tensor might have been defined in terms of the tensor product by setting

$$
\left|x_{1} x_{2}\right|=x_{1} x_{2}-x_{2} x_{1} .
$$

In the same way, if we form the expression ${ }^{4}$

$$
f\left(x_{1}, x_{2}, \cdots, x_{r}\right)=\Sigma \operatorname{sgn}\left(i_{1}, i_{2}, \cdots, i_{r}\right) x_{i_{1}} x_{i_{2}} \cdots x_{i_{r}}
$$

and expand it in terms of the coordinates of the $x$ 's and the fundamental units; it is readily shown that the result is

$$
\sum_{(i)}^{*}\left|\xi_{1 i_{1}} \xi_{2 i_{2}} \cdots \xi_{r i_{r}}\right| f\left(e_{i_{1}}, e_{i_{2}}, \cdots, e_{i_{r}}\right) .
$$

${ }^{4}$ The determinant of a square array of vectors $x_{i j}(i, j=1,2, \cdots, r)$ may be defined as

$$
\left|\dot{x}_{i j}\right|=\Sigma \operatorname{sgn}\left(i_{1}, i_{2}, \cdots, i_{r}\right) x_{1 i_{1}} x_{2 i_{2}} \cdots x_{* i i_{r}} .
$$

In this definition the row marks are kept in normal order and the column marks permuted; a different expression is obtained if the rôles of the row and column marks are interchanged but, as these determinants seem to have little intrinsic interest, it is not worth while to develop a notation for the numerous variants of the definition given above.

Here the scalar multipliers are the same as the coordinates of $\left|x_{1} x_{2} \cdots x_{r}\right|$ and hence the definition of $\$ 1.11$ may now be replaced by

$$
\left|x_{1} x_{2} \cdots x_{r}\right|=\Sigma \operatorname{sgn}\left(i_{1}, i_{2}, \cdots, i_{r}\right) x_{i_{1}} x_{i_{2}} \cdots x_{i_{r}}
$$

which justifies the notation used. We then have

$$
\begin{equation*}
\left|x_{1} x_{2} \cdots x_{r}\right|=\sum_{(i)}^{*}\left|\xi_{1 i_{1}, \xi_{2 i_{2}}} \cdots \xi_{r i_{r}}\right|\left|e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right| . \tag{37}
\end{equation*}
$$

It is easily seen that the tensors $\left|e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right|$ are linearly independent and (37) therefore shows that they form a basis for the set of vectors of grade $r$. Any expression of the form

$$
\Sigma \xi_{i_{1} i_{2}} \cdots i_{r}\left|e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right|
$$

is called a vector of grade $r$ and a vector of the form (37) is called a pure vector of grade $r$.
5.10 The direct product. If $A_{i}=\left\|a_{p q}^{(i)}\right\|(i=1,2, \cdots, r)$ is a sequence of matrices of order $n$, then

$$
\begin{align*}
A_{1} x_{1} A_{2} x_{2} \cdots A_{r} x_{r} & =\sum_{i, j} a_{i_{1} i_{1}}^{(1)} a_{i_{2} i_{2}}^{(2)} \cdots a_{i_{i} j_{r}}^{(r)} \xi_{1 j_{1}} \xi_{2 j_{2}} \cdots \xi_{r i_{r}} e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}  \tag{38}\\
& =\mathfrak{H}\left(x_{1} x_{2} \cdots x_{r}\right)
\end{align*}
$$

where $\mathfrak{A}$ is a linear homogeneous tensor function of $x_{1} x_{2} \cdots x_{r}$, that is, a matrix in space of $n^{r}$ dimensions. This matrix is called the direct product ${ }^{5}$ of $A_{1}, A_{2}, \cdots, A_{r}$ and is denoted by $A_{1} \times A_{2} \times \cdots \times A_{r}$. Obviously

$$
\begin{equation*}
A_{1} B_{1} \times A_{2} B_{2} \times \cdots=\left(A_{1} \times A_{2} \times \cdots\right)\left(B_{1} \times B_{2} \times \cdots\right) \tag{39}
\end{equation*}
$$

and the form of (38) shows that

$$
\begin{equation*}
\left(A_{1} \times A_{2} \times \cdots\right)^{\prime}=A_{1}^{\prime} \times A_{2}^{\prime} \times \cdots \tag{40}
\end{equation*}
$$

From (39) we have, on putting $r=1$ for convenience,'

$$
A_{\perp} \times A_{2} \times A_{3}=\left(A_{1} \times 1 \times 1\right)\left(1 \times A_{2} \times 1\right)\left(1 \times 1 \times A_{3}\right)
$$

Making $A_{i}=1(i=2,3, \cdots, r)$ in (38) we have

$$
A_{1} x_{1} x_{2} \cdots, x_{r}=\Sigma a_{i_{1}, j}^{(1)} \xi_{j_{1}} \xi_{2 i_{2}} \cdots \xi_{r_{i}, e_{i} e_{i} i_{i_{2}}}^{\cdots} e_{i_{r}}
$$

and hence the determinant of the corresponding matrix equals $\left|A_{1}\right|^{n^{r-1}}$. Treating the other factors in the same way we then see that

$$
\begin{equation*}
\left|A_{1} \times A_{2} \times \cdots \times A_{r}\right|=\left|A_{1} A_{2} \cdots A_{r}\right|^{n^{r-1}} \tag{41}
\end{equation*}
$$

Again if as in $\S 5.04$ we take $x_{1}$ as an invariant vector of $A_{1}, x_{2}$ as an invariant vector of $A_{2}$, and so on, and denote the roots of $A_{i}$ by $\lambda_{i j}$, we see that the roots
${ }^{5}$ This definition may be generalized by taking $x_{1}, x_{2}, \cdots$ as vectors in different spaces of possibly different orders. See also §7.03.
of $A_{1} \times A_{2} \times \cdots \times A_{r}$ are the various products $\lambda_{1_{1} i_{1}} \lambda_{2 j_{2}} \cdots \lambda_{r j_{r}}$. When the roots of each matrix are distinct, this gives equation (41) and, since this is an integral relation among the coefficients of the $A$ 's, it follows that it is true in general.

An important particular case arises when each of the matrices in (38) equals the same matrix $A$; the resultant matrix is denoted by $\Pi_{r}(A)$, that is

$$
\begin{equation*}
\Pi_{r}(A)=A \times A \times \cdots \quad(r \text { factors }) \tag{42}
\end{equation*}
$$

It is sometimes called the product transformation. Relations (39), (40), and (41) then become

$$
\begin{equation*}
\Pi_{r}(A B)=\Pi_{r}(A) \Pi_{r}(B), \Pi_{r}(A)^{\prime}=\Pi_{r}\left(A^{\prime}\right),\left|\Pi_{r}(A)\right|=|A|^{r n^{r}-1} \tag{43}
\end{equation*}
$$

5.11 Induced or power matrices. If $x_{1}, x_{2}, \cdots, x_{r}$ are arbitrary vectors, the symmetric expression obtained by forming their products in every possible order and adding is called a permanent. It is usually denoted by ${ }^{\dagger} x_{1} x_{2} \cdots x_{r}{ }^{+}$ but it will be more convenient here to denote it by $\left\{x_{1} x_{2} \cdots x_{r}\right\}$; and similarly, if $\alpha_{i j}$ is a square array of scalars, we shall denote by $\left\{\alpha_{11} \alpha_{22} \cdots \alpha_{r r}\right\}$ the function $\Sigma \alpha_{1 i_{1}} \alpha_{2 i_{8}} \cdots \alpha_{r_{i}}$ in which the summation stretches over every permutation of $1,2, \cdots, r$.

If some of the $x$ 's are equal, the terms of $\left\{x_{1} x_{2} \cdots x_{r}\right\}$ become equal in sets each of which has the same number of terms. If the $x$ 's fall into $s$ groups of $i_{1}, i_{2}, \cdots, i_{s}$ members, respectively, the members in each group being equal to one another, then

$$
\frac{\left\{x_{1} x_{2} \cdots x_{r}\right\}}{i_{1}!i_{2}!\cdots i_{0}!} \quad\left(\Sigma i_{j}=r\right)
$$

has integral coefficients. For the present we shall denote this expression by $\left\{x_{1} x_{2} \cdots x_{r}\right\}^{*}$, but sometimes it will be more convenient to use the more explicit notation

$$
\left\{\begin{array}{cccc}
x_{1} & x_{2} & \cdots & x_{s} \\
i_{1} & i_{2} & \cdots & i_{s}
\end{array}\right\}
$$

in which $i_{1}$ of the $x$ 's equal $x_{1}, i_{2}$ equal $x_{2}$, etc.; this notation is, in fact, that already used in $\S 2.08$, for instance,

$$
\begin{aligned}
& \left\{\begin{array}{lll}
x & x & y
\end{array}\right\}=2 x^{2} y+2 x y x .+2 y x^{2} \\
& \left\{\begin{array}{ll}
x & y \\
2 & 1
\end{array}\right\}=x^{2} y+x y x+y x^{2}=\frac{1}{2}\{x x y\}
\end{aligned}
$$

The same convention applies immediately to $\left\{\alpha_{11} \alpha_{2 ?} \cdots \alpha_{r r}\right\}$.
In the notation just explained we have

$$
\begin{equation*}
\left\{x_{1} x_{2} \cdots x_{r}\right\}=\Sigma^{\prime \prime}\left\{\xi_{1_{i} i_{1}} z_{2 i_{2}} \cdots \xi_{r i_{r}}\right\}^{*}\left\{e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right\} \tag{44}
\end{equation*}
$$

where the summation $\Sigma^{\prime \prime}$ extends over all combinations $i_{1} i_{2} \cdots i_{r}$ of the first $n$ integers repetition being allowed. This shows that the set of all permanents
of grade $r$ has the basis $\left\{e_{i_{1}} e_{i_{2}} \cdots e_{i_{r}}\right\}$ of order $(n+r-1)!/ r!(n-1)!$. From (44) we readily derive

$$
\begin{equation*}
\left\{A x_{1} A x_{2} \cdots A x_{r}\right\}=\sum_{i, j}^{\prime \prime}\left\{a_{i_{1 i} i_{1}} a_{i_{2} i_{2}} \cdots a_{i r r_{r}}\right\}^{*}\left\{\xi_{1_{i_{1}}} \cdots \xi_{r j_{r}}\right\}^{*}\left\{e_{i_{1}} \cdots e_{i_{r}}\right\} \tag{45}
\end{equation*}
$$

which is a linear tensor form in $\left\{x_{1} x_{2} \cdots x_{r}\right\}$. We may therefore set

$$
\begin{equation*}
\left\{A x_{1} A x_{2} \cdots A x_{r}\right\}=P_{r}(A)\left\{x_{1} x_{2} \cdots x_{r}\right\} \tag{46}
\end{equation*}
$$

where $P_{r}(A)$ is a matrix of order $(n+r-1)!/ r!(n-1)!$ whose coordinates are the polynomials in the coordinates of $A$ which are given in (45); this matrix is called the $r$ th induced or power matrix of $A$. As with $C_{r}(A)$ and $\Pi_{r}(A)$ it follows that

$$
\begin{gather*}
P_{r}(A B)=P_{r}(A) P_{r}(B), P_{r}(A)^{\prime}=P_{r}\left(A^{\prime}\right), \\
\left|P_{r}(A)\right|=|A|^{\binom{n+r-1}{r-1}} ; \tag{47}
\end{gather*}
$$

also the roots of $P_{r}(A)$ are the various products of the form $\lambda_{1}{ }^{\alpha_{1}} \lambda_{2}{ }^{\alpha_{2}} \cdots \lambda_{r}{ }^{\alpha}{ }^{\prime}$ for which $\Sigma \alpha_{i}=r$.
5.12 Associated matrices. The matrices considered in the preceding sections have certain common properties; the coordinates of each are functions of the variable matrix $A$ and, if $T(A)$ stands for any one of them, then

$$
\begin{equation*}
T(A B)=T(A) T(B) \tag{48}
\end{equation*}
$$

Following Schur, who first treated the general problem of determining all such matrices, we shall call any matrix with these properties an associated matrix. If $S$ is any constant matrix in the same space as $T(A)$, then $T_{1}(A)=S T(A) S^{-1}$ is clearly also an associated matrix; associated matrices related in this manner are said to be equivalent.

Let the orders of $A$ and $T(A)$ be $n$ and $m$ respectively and denote the corresponding identity matrices by $1_{n}$ and $1_{m}$; then from (48)

$$
\begin{equation*}
T^{2}\left(1_{n}\right)=T\left(1_{n}\right), \quad T\left(1_{n}\right) T(A)=T(A)=T \dot{(A)} T\left(1_{n}\right) \tag{49}
\end{equation*}
$$

If $s$ is the rank of $T\left(1_{n}\right)$, we can find a matrix $S$ which transforms $T\left(1_{n}\right)$ into a diagonal matrix with $s$ 1's in the main diagonal and zeros elsewhere; and we may without real loss of generality assume that $T\left(1_{n}\right)$ has this form to start with, and write

$$
T\left(1_{n}\right)=\left\|\begin{array}{ll}
1_{s} & 0 \\
0 & 0
\end{array}\right\|
$$

The second equation of (49) then shows that $T(A)$ has the form

$$
T(A)=\left\|\begin{array}{cc}
T_{s}(A) & 0 \\
0 & 0
\end{array}\right\|
$$

and we shall therefore assume that $s=m$ so that $T\left(1_{n}\right)=1_{m}$. It follows from this that $|T(A)| \not \equiv 0$ so that $T(A)$ is not singular for every $A$; we shall then say that $T$ is non-singular.

A non-singular associated matrix $T(A)$ is reducible (cf. §3.10) if it can be expressed in the form $T(A)=T_{1}(A)+T_{2}(A)$ where, if $E_{1}=T_{1}\left(1_{n}\right), E_{2}=$ $T_{2}\left(1_{n}\right)$, so that $E_{1}+E_{2}=1_{m}$, then

$$
\begin{gathered}
T_{1}(A)=E_{1} T(A) E_{1}, \quad T_{2}(A)=E_{2} T(A) E_{2} \\
E_{1} T(A) E_{2}=0=E_{2} T(A) E_{1}
\end{gathered}
$$

so that

$$
\begin{gathered}
E_{1}^{2}=E_{1}, E_{2}^{2}=E_{2} \\
E_{1} E_{2}=0=E_{2} E_{1}
\end{gathered}
$$

and there is therefore an equivalent associated matrix $t(A)$ which has the form

$$
t(A)=\left\|\begin{array}{cc}
t_{1}(A) & 0 \\
0 & t_{2}(A)
\end{array}\right\|
$$

When $T(A)$ is reducible in this manner we have

$$
\begin{aligned}
T_{1}(A B) & =E_{1} T(A B) E_{1}=E_{1} T(A) T(B) E_{1} \\
& =E_{1} T(A)\left(E_{1}+E_{2}\right) T(B) E_{1} \\
& =E_{1} T(A) E_{1} T(B) E_{1}=T_{1}(A) T_{1}(B)
\end{aligned}
$$

so that $T_{1}(A)$ and $T_{2}(A)$ are separately associated matrices. We may therefore assume $T(A)$ irreducible without loss of generality since reducible associated matrices may be built up out of irreducible ones by reversing the process used above.
5.13 We shall now show that, if $\lambda$ is a scalar variable, then $T(\lambda)$ is a power of $\lambda$. To begin with we shall assume that the coordinates of $T(\lambda)$ are rational functions in $\lambda$ and that $T(1)$ is finite; we can then set $T(\lambda)=T_{1}(\lambda) / f(\lambda)$ where $f(\lambda)$ is a scalar polynomial whose leading coefficient is 1 and the coordinates of $T_{1}(\lambda)$ are polynomials whose highest common factor has no factor in common with $f(\lambda)$. If $\mu$ is a second scalar variable, (48) then gives

$$
\frac{T_{1}(\lambda) T_{1}(\mu)}{f(\lambda) f(\mu)}=\frac{T_{1}(\lambda \mu)}{f(\lambda \mu)}
$$

hence $f(\lambda \mu)$ is a factor of $f(\lambda) f(\mu)$, from which it follows readily that $f(\lambda \mu)$ $=f(\lambda) f(\mu)$; so that $f(\lambda)$ is a power of $\lambda$ and also

$$
\begin{equation*}
T_{1}(\lambda \mu)=T_{1}(\lambda) T_{1}^{\prime}(\mu) \tag{50}
\end{equation*}
$$

We also have $f(1)=1$ and hence $T_{1}\left(1_{n}\right)=T\left(1_{n}\right)=1_{m}$.
Let $T_{1}(\lambda)=F_{0}+\lambda F_{1}+\cdots+\lambda^{s} F_{s}\left(F_{s} \neq 0\right)$; then from (50)

$$
F_{0}+\lambda \mu F_{1}+\cdots+\lambda^{\varepsilon} \mu^{\bullet} F_{s}=\left(F_{0}+\lambda F_{1}+\cdots\right)\left(F_{0}+\mu F_{1}+\cdots\right)
$$

which gives

$$
F_{i}^{2}=F_{i}, F_{i}^{\prime} F_{j}=0(i \neq j), \quad(i, j=0,1, \cdots, s)
$$

Now

$$
T_{1}(\lambda) T(A)=f(\lambda) T(\lambda) T(A)=f(\lambda) T(\lambda A)=T(A) T_{1}(\lambda)
$$

therefore

$$
\Sigma F_{i} T(A) \lambda^{i}=\Sigma T(A) F_{i} \lambda^{i}
$$

and hence on comparing powers of $\lambda$ we have

$$
F_{i} T(A)=T(A) F_{i}
$$

and, since $\Sigma F_{i}=T_{1}(1)=1_{m}$ and we have assumed that $T(A)$ is irreducible, it follows that every $F_{i}=0$ except $F_{s}$, which therefore equals $1_{m}$. Hence $T_{1}(\lambda)=\lambda^{\circ}$ and, since $f(\lambda)$ is a power of $\lambda$, we may set

$$
\begin{equation*}
T(\lambda)=\lambda^{r} \tag{51}
\end{equation*}
$$

Since $T(\lambda A)=T(\lambda) T(A)=\lambda^{r} T(A)$, we have the following theorem.
Theorem 9. If $T(A)$ is irreducible, and if $T(\lambda)$ is a rational function of the scalar variable $\lambda$, then $T(\lambda)=\lambda^{r}$ and the coordinates of $T(A)$ are homogeneous functions of order $r$ in the coordinates of $A$.

The restriction that $T(\lambda)$ is rational in $\lambda$ is not wholly necessary. For instance, if $q$ is any whole number and $\epsilon$ a corresponding primitive root of 1 , then $T^{q}(\epsilon)=1_{m}$ and from this it follows without much difficulty that $T(\epsilon)=\epsilon^{\circ}$ where $s$ is an integer which may be taken to be the same for any finite number of values of $q$. It follows then that, if $T(\lambda)=\left\|t_{i j}(\lambda)\right\|$, the functions $t_{i j}(\lambda)$ satisfy the equation

$$
t_{i j}(\epsilon \lambda)=\epsilon^{8} t_{i j}(\lambda)
$$

and under very wide assumptions as to the nature of the functions $t_{i j}$ it follows from this that $T(\lambda)$ has the form $\lambda^{r}$. Again, if we assume that $T(\lambda)=\lambda^{\alpha} \sum_{-\alpha}^{\alpha} T_{r} \lambda^{r}$, then $T(\lambda) T(\mu)=T(\lambda \mu)$ gives immediately

$$
T_{r \mu^{r}+\alpha}=T(\mu)
$$

so that only one value of $r$ is admissible and for this value $T_{r}=1$ as before.
5.14 If the coordinates of $T(A)$ are rational functions of the coordinates $a_{i j}$ of $A$, so that $r$ is an integer, we can set $T(A)=T_{1}(A) / f(A)$ where the coordinates of $T_{1}(A)$ are integral in the $a_{i j}$ and $f(A)$ is a scalar polynomial in these variables which has no factor common to all the coordinates of $T_{1}(A)$. As in (50) we then have

$$
T_{1}(A B)=T_{1}(A) T_{1}(B), f(A B)=f(A) f(B)
$$

It follows from the theory of scalar invariants that $f(A)$ can be taken as a positive integral power of $|A|$; we shall therefore from this point on assume that the coordinates of $T(A)$ are homogeneous polynomials in the coordinates of $A$ unless the contrary is stated explicitly. We shall call $r$ the index of $T(A)$.

Theorem 10. If $T(A)$ is an associaied matrix of order $m$ and index $\dot{r}$, and if the roots of $A$ are $\alpha_{1}, \alpha_{2}, \cdots, \alpha_{n}$, then the roots of $T(A)$ have the form $\alpha_{1}^{r_{1}} \alpha_{2}^{r_{2}} \cdots$ $\cdot \alpha_{n}^{r_{n}}$ where $\Sigma r_{i}=r$. The actual choice of the exponents $r$ depends on the particular associated $\dot{m} a t r i x$ in question but, if $\alpha_{1}^{r_{1}} \alpha_{2}^{r_{2}} \cdots \alpha_{n}^{r_{n}}$ is one root, all the distinct quantities obtained from it by permuting the $\alpha$ 's are also roots.

If the roots of $A$ are arbitrary variables; then $A$ is similar to a diagonal $\operatorname{matrix} A_{1}=\Sigma \alpha_{i} e_{i i}$. We can express $T\left(A_{1}\right)$ as a polynomial ${ }^{6}$ in the $\alpha$ 's, say

$$
\begin{equation*}
T\left(A_{1}\right)=\Sigma \alpha_{1}^{r_{1}} \alpha_{2}^{r_{2}} \cdots \alpha_{n}^{r} F_{r_{1} r_{2}} \cdots r_{n} \tag{52}
\end{equation*}
$$

where the $F$ 's are constant matrices. If now $B=\Sigma \beta_{i} e_{i i}$ is a second variable diagonal matrix, the relation $T\left(A_{1} B\right)=T\left(A_{1}\right) T(B)$ gives as in (50)

$$
\begin{align*}
F_{r_{1} r_{2} \cdots r_{n}}^{2} & \left.\left.=F_{r_{1} r_{2} \cdots r_{n}}, \ldots\right) \neq\left(s_{1}, s_{2} \cdots\right)\right) \\
F_{r_{1} r_{2} \cdots r_{n}} F_{s_{1} s_{2}} \cdots s_{n} & =0\left(\left(r_{1}, r_{2}, \cdots\right)\right. \tag{53}
\end{align*}
$$

and hence $T\left(A_{1}\right)$ can be expressed as a diagonal matrix with roots of the required form; these roots may of course be multiple since the rank of $F_{r_{1} \ldots r_{n}}$ is not necessarily 1 , the elementary divisors are, however, simple.
Since the associated matrices of similar matrices are similar, it follows that the roots of the characteristic equation of $T(A)$ are given by those terms in (52) for which $F_{r_{1} r_{2} \ldots r_{n}} \neq 0$; and, since this equation has coefficients which are polynomials in the coordinates of $A$, the roots of $T(A)$ remain in this form even when the roots of $A$ are not necessarily all different.

The rest of the theorem follows from the fact that the trace of $T\left(A_{1}\right)$ equals that of $T(A)$ which is rational in the coordinates of $A$ and is therefore symmetric in the $\alpha$ 's.

Theorem 11. The value of the determinant of $T(A)$ is $|A|^{\text {rm/n }}$ and $r m / n$ is an integer.

For $T(A) T(\operatorname{adj} A)=T(|A|)=|A|^{\circ}$ and therefore $|T(A)|$ is a power of $|A|$, say $|A|^{\circ}$. But $T(A)$ is a matrix of order $m$ whose coordinates are polynomials in the coordinates of $A$. Hence $s n=m r$ and $r m / n$ is an integer.
5.15 Transformable systems. From a scalar point of view each of the associated matrices discussed in $\S \S 5.03-5.11$ can be characterized by a set of scalar functions $f_{k}(k=1,2, \cdots, m)$ of one or more sets of variables ( $\xi_{i}$,
${ }^{6}$ If we merely assume that $T\left(A_{1}\right)$ is a convergent series of the form (52), equation (53) still holds. It follows that there are only a finite number of terms in (52) since (53) shows that there is no linear relation among those $F_{r_{1}} \ldots r_{n}$ which are not zero. Let $F_{i}$ be the sum of those $F_{r_{1}} \ldots r_{n}$ for which $\Sigma r_{i}$ has a fixed value $\rho_{i}$; then $T(\lambda)=\Sigma \lambda^{\rho_{i}} F_{i}$, and as before only one value of $\rho_{1}$ is admissible when $T(A)$ is irreducible.
$j=1,2, \cdots, n), i=1,2, \cdots, r$, which have the following property: if the $\xi$ 's are subjected to a linear transformation

$$
\xi_{i j}^{\prime}=\sum_{s=1}^{n} a_{j s} \xi_{i s} \quad(j=1,2, \cdots, n ; i=1,2, \cdots, r)
$$

and if $f_{k}^{\prime}$ is the result of replacing $\xi_{i j}$ by $\xi_{i j}^{\prime}$ in $f_{k}$, then

$$
\mathbf{f}_{k}^{\prime}=\sum_{s=1}^{m} \alpha_{k s} f
$$

where the $\alpha$ 's are functions of the $a_{i j}$ and are independent of the $\xi$ 's. For instance, corresponding to $C_{2}(A)$ we have

$$
f_{i} \equiv f_{p q}=\left|\begin{array}{ll}
\xi_{1 p} & \xi_{1 q} \\
\xi_{2 p} & \xi_{2 q}
\end{array}\right| \quad(p, q=1,2, \cdots, n ; p<g)
$$

for which

$$
\alpha_{i j} \equiv \alpha_{p q, r s}=\left|\begin{array}{cc}
a_{p r} & a_{q r} \\
a_{p s} & a_{q s}
\end{array}\right|
$$

We may, and will, always assume that there are no constant multipliers such that $\Sigma \lambda_{i} f_{i}=0$. Such systems of, functions were first considered by Sylvester; they are now generally called transformable systems.

If we put $T(A)=\left\|\alpha_{i j}\right\|$, we have immediately $T(A B)=T(A) T(B)$, and consequently there is an associated matrix corresponding to every transformable system. Conversely, there is a transformable system corresponding to an associated matrix. For if $X=\left\|\xi_{i j}\right\|$ is a variable matrix and $c$ an arbitrary constant vector in the space of $T(A)$, then the coordinates of $T(X) c$ form a transformable system since $T(A) T(X) c=T(A X) c$ and $c$ can be so determined that there is no constant vector $b$ such that $S b T(X) c \equiv 0$.

The basis $f_{k}(k=1,2, \cdots, m)$ may of ${ }^{\cdot}$ course be replaced by any basis which is equivalent in the sense of linear dependence, the result of such a change being to replace $T(A)$ by an equivalent associated matrix. If in particular there exists a basis

$$
g_{1}, g_{2}, \cdots, g_{k_{1}}, h_{1}, h_{2}, \cdots, h_{k_{2}} \quad\left(k_{1}+k_{2}=k\right)
$$

such that the $g$ 's and the $h$ 's form separate transformable systems, then $T(A)$ is reducible; and conversely, if $T(A)$ is reducible, there always exists a basis of this kind.
5.16 Transformable linear sets. If we adopt the tensor point of view rather than the scalar one, an associated matrix is found to be connected with a linear set $\mathfrak{F}$ of constant tensors, derived from the fundamental units $e_{i}$, such that, when $e_{i}$ is replaced by $A e_{i}(i=1,2, \cdots, n)$ in the members of the basis of $\mathfrak{F}$, then the new tensors are linearly dependent on the old; in other words
the set $\mathfrak{F}$ is invariant as a whole under any linear transformation $A$ of the fundamental units. For instance, in the case of $C_{2}(A)$ cited above, $\mathfrak{F}$ is the linear set defined by

$$
\left|e_{i} e_{j}\right| \quad(i, j=1,2, \cdots, n ; i<j)
$$

We shall call a set which has this property a transformable linear set.
Let $u_{1}, u_{2}, \cdots, u_{m}$ be a transformable linear set of tensors of grade $r$ and let $u_{i}^{\prime}$ be the tensor that results when $e_{j}$ is replaced by $A e_{j}(j=1,2, \cdots, n)$ in $u_{i}$. Since the set is transformable, we have

$$
u_{i}^{\prime}=\sum_{j} \alpha_{j i} u_{j}=T(A) u_{i} \quad(i=1,2, \cdots, m)
$$

where the $\alpha_{i j}$ are homogeneous polynomials in the coordinates of $A$ of degree $r$. If we employ a second transformation $B$, we then have

$$
u_{i}^{\prime \prime}=T(A) T(B) u_{i}, \quad u_{i}^{\prime \prime}=T(A B) u_{i} \quad(i=1,2, \cdots, m)
$$

and therefore $T(A)$ is an associated matrix.
We have now to show that there is a transformable linear set corresponding to every associated matrix. In doing this it is convenient to extend the notation Suv to the case where $u$ and $v$ are tensors of grade $r$. Let $E_{i}\left(i=1,2, \cdots, n^{r}\right)$ be the unit tensors of grade $r$ and

$$
u=\Sigma \psi_{i} E_{i}, v=\Sigma \varphi_{i} E_{i}
$$

any tensors of grade $r$; we then define $S u v$ by

$$
S u v=\left(\sum_{1}^{n^{r}} \psi_{i} \varphi_{i}\right) / r!
$$

where the numerical divisor is introduced solely in order not to disagree with the definition of $\S 5.02$.

Let $x_{i}=\Sigma \xi_{i j} e_{j}(i=1,2, \cdots, r)$ be a set of variable vectors and $X_{i}(i=1$, $\cdot 2, \cdots, s)$ the set of tensors of the form $x_{1}^{j 1} x_{2}^{j 2} \cdots x_{r}^{j r}\left(\Sigma j_{i}=r\right)$; we can then put any product $\xi_{11}^{\beta_{11}} \xi_{12}^{\beta_{12}} \cdots \xi_{r n}^{\beta_{i n}}$ for which $\Sigma \beta_{i j}=r$ in the form $k S E_{i} X_{i}, k$ being a numerical factor. This can be done in more than one way as a rule; in fact, if $\sum_{j} \beta_{i j}=\beta_{i}$, then

$$
\xi_{11}^{\beta_{11}} \cdots \xi_{1 n}^{\beta_{1 n}}=\frac{1}{\beta_{1}!} S e_{1}^{\beta_{11}} \cdots e_{n}^{\beta_{1 n}} x_{1}^{\beta_{1}}
$$

and from the definition of $S u v$ it is clear that the factors in $e_{1}^{\beta_{11}} \cdots e_{n}^{\beta_{1 n}}$ can be permuted in any way without altering the value of the scalar. It follows that

$$
\xi_{11}^{\beta_{11}} \cdots \xi_{1 n}^{\beta_{1 n}}=\frac{1}{\beta_{1!}!\beta_{12}!\cdots \beta_{1 n}!} S\left\{\begin{array}{ccc}
e_{1} & \cdots & e_{n} \\
\beta_{11} & \cdots & \beta_{1 n}
\end{array}\right\} x_{1}^{\beta_{1}}
$$

and repeating this process we get

$$
k_{1} \xi_{11}^{\beta_{11}} \xi_{12}^{\beta_{12}} \cdots \xi_{r n}^{\beta_{r n}}=S\left\{\begin{array}{ccc}
e_{1} & \cdots & e_{n} \\
\beta_{11} & \cdots & \beta_{1 n}
\end{array}\right\} \cdots\left\{\begin{array}{ll}
\cdots & e_{n} \\
\cdots & \beta_{r n}
\end{array}\right\} x_{1}^{\beta_{1}} x_{2}^{\beta_{3}} \cdots \cdots x_{r}^{\beta_{r}}
$$

where $k_{1}$ is a numerical factor whose value is immaterial for our present purposes.
If $f$ is any homogeneous polynomial in the variables $\xi_{i j}$ of degree $\rho$, it can be expressed uniquely in the form

$$
f=\Sigma \Sigma \varphi_{\beta_{11}} \cdots \beta_{r n} S\left\{\begin{array}{lll}
e_{1} & \cdots & e_{n} \\
\beta_{11} & \cdots & \beta_{1 n}
\end{array}\right\} \cdots\left\{\begin{array}{ccc}
e_{1} & \cdots & e_{n} \\
\beta_{r 1} & \cdots & \beta_{r n}
\end{array}\right\} x_{1}^{\beta_{1}} \cdots \cdots x_{r}^{\beta_{r}}
$$

where the inner summation extends over the partitions of $\beta_{i}$ into $\beta_{i 1}, \beta_{i 2}, \cdots, \beta_{i n}$ ( $i=1,2, \cdots, r$ ) and the outer over all values of $\beta_{1}, \beta_{2}, \cdots, \beta_{r}$ for which $\Sigma \beta_{i}=\rho$. We may therefore write

$$
f=\sum_{1}^{s} S F_{j} X_{j}
$$

where, as above, $X_{i}=x_{1}^{\beta_{1}} x_{2}^{\beta_{2}} \cdots x_{r}^{\beta_{r}}$ and

$$
F_{j} \equiv F_{\beta_{1} \beta_{2}} \ldots \beta_{r}=\Sigma \varphi_{\beta_{11}} \ldots \beta_{r n}\left\{\begin{array}{lll}
e_{1} & \cdots & e_{n} \\
\beta_{11} & \cdots & \beta_{1 n}
\end{array}\right\} \cdots\left\{\begin{array}{lll}
e_{1} & \cdots & e_{n} \\
\beta_{r 1} & \cdots & \beta_{r n}
\end{array}\right\} .
$$

The expression of $f$ in this form is unique. In the first place, $F_{i} \neq 0$ unless each $\varphi_{\beta_{11}} \ldots \beta_{r n}$ is zero, since the set of tensors of the form

$$
\left\{\begin{array}{ccc}
e_{1} & \cdots & e_{n} \\
\beta_{11} & \cdots & \beta_{1 n}
\end{array}\right\} \cdots\left\{\begin{array}{lll}
e_{1} & \cdots & e_{n} \\
\beta_{r 1} & \cdots & \beta_{r n}
\end{array}\right\} \quad\left(\Sigma \beta_{i j}=\rho\right)
$$

are clearly linearly independent. Further, if $\Sigma S F_{j} X_{j} \equiv 0$, then each $S F_{j} X_{i}$ is zero since each gives rise to terms of different type in the $\xi_{i j}$; and finally the form of $F_{j}$ shows that $S F_{j} X_{j}=0$ only if $F_{i}=0$ since in

$$
S F_{i} X_{i}=k_{1} \Sigma \varphi_{\beta_{11}} \cdots \beta_{r n} \xi_{11}^{\xi_{1}^{11}} \cdots \xi_{r n}^{\beta_{r n}}
$$

each term of the summation is of different type in the $\xi_{i j}$.
Let $\left(f_{k}\right)$ be a transformable system; we can now write uniquely

$$
\begin{equation*}
f_{k}=\sum_{i} S F_{k j} X_{j}(k=1,2, \cdots, m) \tag{54}
\end{equation*}
$$

and we may set

$$
F=\sum_{i}^{n^{r}} f_{i} E_{i}=\sum_{i, j} E_{i} S F_{i j} X_{i}
$$

where $f_{i} \equiv 0$ when $i>m$. If we transform the $x$ 's by $A=\left\|a_{i j}\right\|$ and denote $\Pi_{r}(A)$ temporarily by $\Pi$, then $X_{i}$, becomes $\Pi X_{j}$ and $F$ is transformed into $F^{*}$ where

$$
\begin{equation*}
F^{*}=\sum_{i, j} E_{i} S F_{i j} \Pi X_{i}=\sum_{i, j} E_{i} S \Pi^{\prime} F_{i j} \cdot X_{i} . \tag{55}
\end{equation*}
$$

But the $f$ 's form a transformable system and hence by this transformation $f_{i}$ becomes

$$
f_{i}^{\prime}=\sum_{k} \alpha_{i k} f_{k}
$$

so that

$$
\begin{equation*}
F^{*}=\sum_{k, i} \alpha_{i k} f_{k} E_{i}=\sum_{i} E_{i} S \sum_{k} \alpha_{i k} \sum_{i} F_{k j} X_{j} \tag{56}
\end{equation*}
$$

Comparing (55) and (56) we have

$$
\begin{equation*}
\sum_{j} S\left[\sum_{k} \alpha_{i k} F_{k j}-\Pi^{\prime} F_{i j}\right] X_{i}=0 \tag{57}
\end{equation*}
$$

and therefore, as was proved above, each of the terms of the summation is zero, that is,

$$
\begin{equation*}
\Pi^{\prime} F_{i j}=\sum_{k} \alpha_{i k} F_{k j} \tag{58}
\end{equation*}
$$

and therefore, if $j$ is kept fixed, the linear set

$$
\begin{equation*}
\left(F_{1 j}, F_{2 j}, \cdots\right) \tag{59}
\end{equation*}
$$

is transformable provided $F_{1 j}, F_{2 j}, \cdots$ are linearly independent.
If there is no $j$ for which the set (59) is linearly independent we proceed as follows. Let $f_{i j}=S F_{i j} X_{i}$ so that

$$
\begin{gather*}
f_{1}=f_{11}+f_{12}+\cdots+f_{1 s} \\
f_{2}=f_{21}+f_{22}+\cdots+f_{28}  \tag{60}\\
\cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \\
\cdots \cdots \cdots \cdots \cdots \cdots \cdots \\
f_{m}=f_{m 1}+f_{m 2}+\cdots+f_{m s}
\end{gather*}
$$

If the removal of any column of this array leaves the new $f_{i}$ so defined linearly independent, they form a transformable system which defines the same associated matrix as the original system; we shall therefore suppose that the removal of any column leads to linear relations among the rows, the coefficients of these relations being constants. Remove now the first column; then by non-singular constant combinations of the rows we can make certain of them, say the first $m_{1}$, equal 0 , the remainder being linearly independent. On applying the same transformation to the rows of (60), which leaves it still a transformable system, we see that we may replace (60) by an array of the form

$$
\begin{align*}
& f_{1}=f_{11} \\
& \ldots \ldots \ldots \cdots  \tag{61}\\
& f_{m_{1}}=f_{m_{1} 1} \\
& f_{m_{1}+1}=f_{m_{1}+1,1}+f_{m_{1}+1,2}+\cdots+f_{m_{1}+1, s} \\
& \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \\
& f_{m}=f_{m 1}+f_{m 2}+\cdots \cdots+f_{m s}
\end{align*}
$$

where $f_{m_{1}+i}-f_{m_{1}+i, j}\left(i=1,2, \cdots, m-m_{1}\right)$ are linearly independent. It follows that $f_{1}, \cdots, f_{m_{1}}$ are transformed among themselves and so form a transformable system. For these functions are transformed in the same way as $f_{11}, f_{21} \cdots, f_{m_{1} 1}$, and if the last $m-m_{1}$ rows of (61) were involved in the transformation, this would mean that $f_{11}, \cdots, f_{m_{1}}$, when transformed, would depend on $f_{m_{1}+1, i}$ etc., which is impossible owing to the linear independence of $f_{m_{1}+i}-f_{m_{1}+i, 1}\left(i=1,2, \cdots, m-m_{1}\right)$.

Corresponding to the first column of (61) we have tensors $F_{11}, F_{21}, \cdots, F_{m 1}$ and we may suppose this basis so chosen that $F_{i 1}(i=1,2, \cdots, p)$ are linearly independent and $F_{j 1}=0$ for $j>p$; and this can be done without disturbing the general form of (61). If $p=m$, we have a transformable system of the type we wish to obtain and we shall therefore assume that $p<m$. We may also suppose the basis so chosen that $S \bar{F}_{i 1} F_{i 1}=\delta_{i j}(i, j=1,2, \cdots, p)$ as in Lemma 2, §1.09. It follows from what we have proved above that $F_{11}, F_{21}$, $\cdots, F_{m_{1} 1}$ is a transformable set.

Let $A$ be a real matrix, the corresponding transformation of the $F$ 's being, as in (58),

$$
\begin{equation*}
F_{i 1}^{*}=\sum_{j} \alpha_{i j} F_{j_{1}}=\Pi^{\prime} F_{i 1}, \quad(i=1,2, \cdots, p) \tag{62}
\end{equation*}
$$

we then have

$$
\begin{equation*}
\bar{F}_{i 1}^{*}=\sum_{j} \bar{\alpha}_{i j} \bar{F}_{j 1}=\Pi^{\prime}(A) \bar{F}_{i 1} \tag{63}
\end{equation*}
$$

so that the $\bar{F}_{i 1}$ also forms a transformable set. Since $F_{11}, \cdots, F_{m_{1} 1}$ form a transformable set, $\alpha_{i j}$ and $\bar{\alpha}_{i j}$ are 0 when $i>m_{1}$ and $j \leq m_{1}$ no matter what matrix $A$ is. Now

$$
\alpha_{i j}=S \bar{F}_{j_{1}} F_{i 1}^{*}=S \bar{F}_{j 1} \Pi^{\prime}(A) F_{i 1}=S \Pi(A) \bar{F}_{j_{1}} F_{i_{1}}=S \Pi^{\prime}\left(A^{\prime}\right) \bar{F}_{j_{1}} F_{i_{1}}
$$

which equals 0 for $i \leq m_{1}, j>m_{1}$ since by (63) $\Pi^{\prime}\left(A^{\prime}\right) \bar{F}_{j 1}$ is derived from $\bar{F}_{j 1}$ by the transformation $A^{\prime}$ on the $x$ 's and for $j \leq m_{1}$ is therefore linearly dependent on $\bar{F}_{j 1}\left(j=1,2, \cdots, m_{1}\right)$. Hence the last $m-m_{1}$ rows in (61) also form a transformable system, which is only possible if the system $f_{1}, f_{2}, \cdots, f_{m}$ is reducible. If $T(A)$ is irreducible, the corresponding transformable system is irreducible and it follows now that there also corresponds to it an irreducible transformable set of tensors.
5.17 We have now shown that to every associated matrix $T(A)$ of index $r$ and order $m$ there corresponds a transformable linear set of constant tensors $F_{1}, F_{2}, \cdots, F_{m}$ of grade $r$ whose law of transformation is given by (62). Also since $\Pi^{\prime}(A)=\Pi\left(A^{\prime}\right)$, we have

$$
\begin{equation*}
\Pi F_{i}=\Sigma \alpha_{i k}^{\prime} F_{k}, \quad \Pi \bar{F}_{i}=\Sigma \bar{\alpha}_{i k}^{\prime} \bar{F}_{k} \tag{64}
\end{equation*}
$$

where $T\left(A^{\prime}\right)=\left\|\alpha_{i j}^{\prime}\right\|$.

Since $F_{1}, F_{2}, \cdots, F_{m}$ are linearly independent, we can find a supplement to this set in the set of all tensors of grade $r$, say

$$
G_{1}, G_{2}, \cdots, G_{\mu} \quad\left(\mu=n^{r}-m\right)
$$

such that

$$
\begin{equation*}
S \bar{F}_{i} G_{j}=0 \tag{65}
\end{equation*}
$$

It is convenient also to choose bases for both sets such that

$$
\begin{equation*}
S \bar{F}_{i} F_{i}=\delta_{i i}=S \bar{G}_{i} G_{j} \tag{65'}
\end{equation*}
$$

Since the two sets together form a basis for the space of $\Pi$, we can set

$$
\Pi^{\prime} G_{j}=\Sigma \beta_{k j} F_{k}+\Sigma \gamma_{k_{j}} G_{k}
$$

and this gives

$$
\beta_{i j}=S \bar{F}_{i} \Pi^{\prime} G_{j}=S G_{i} \Pi \bar{F}_{i}
$$

which is 0 from (64) and (65), hence the $G$ 's are transformed among themselves by $\Pi^{\prime}$. This means, however, that $\Pi^{\prime}$ is reducible, and when it is expressed in terms of the basis ( $F_{1}, \cdots, F_{m}, G_{1}, \cdots, G_{\mu}$ ), the part corresponding to ( $F_{1}, \cdots, F_{m}$ ) has the form $\left\|\alpha_{i j}\right\|$ and is therefore similar to $T(A)$. Hence:

Theorem 12. Every irreducible associated matrix $T(A)$ of index $r$ is equivalent to an irreducible part of $\Pi_{r}(A)$, and conversely.
5.18 Irreducible transformable sets. If $F$ is a member of a transformable linear set $\mathfrak{F}=\left(F_{1}, F_{2}, \cdots, F_{m}\right)$, the total set of tensors derived from $F$ by all linear transformations of the fundamental units clearly form a transformable linear set which is contained in $\mathfrak{F}$, say $\mathfrak{F}_{1}$; and we may suppose the basis of $\mathfrak{F}$ so chosen that $\mathfrak{F}_{1}=\left(F_{1}, F_{2}, \cdots, F_{k}\right)$ and $S \bar{F}_{i} F_{i}=\delta_{i j}(i, j=1,2, \cdots, m)$. Let $G$ be an element of $\left(F_{k+1}, \cdots, F_{m}\right)$ and $G^{\prime}$ a transform of $G$ so that

$$
G^{\prime}=\sum_{i=1}^{m} \gamma_{i} F_{i}
$$

Then $S \bar{F}_{i} G^{\prime}=\gamma_{i}$. But $S \bar{F}_{i} G^{\prime}=S \bar{F}_{i}^{\prime} G$, where $F_{i}^{\prime}$ is the transform of $F_{i}$ obtained by the transverse of the transformation which produced $G^{\prime}$ from $G$ sothat $\bar{F}_{i}^{\prime}$ is in $\mathfrak{F}_{1}$ for $i \leq k$. Hence $\gamma_{i}=0$ for $i=1,2 ; \cdots, k$, that is, $\left(F_{k+1}, \cdots, F_{m}\right)$ is also a transformable set; and so, when the original set is irreducible, we must have $\mathfrak{F}_{1}=\mathfrak{F}$. If we say that $F$ generates $\mathfrak{F}$, this result may be stated as follows.

Lemma 5. An irreducible transformable linear set is generated by any one of its members.

We may choose $F$ so that it is homogeneous in each $e_{i}$; for if we replace, say, $e_{1}$ by $\lambda \theta_{1}$, then $F$ has the form $\Sigma \lambda^{k} H_{k}$ and by the same argument as in $\S 5.13$, any $H_{k}$ which is not 0 is homogeneous in $e_{1}$ and belongs to $\mathfrak{F}$. A repetition of
this argument shows that we may choose $F$ to be homogeneous in each of the fundamental units which occur in it. If $r$ is the grade of $F$, we may assume that $F$ depends'on $e_{1}, e_{2}, \cdots, e_{s}$, and; if $k_{1}, k_{2}, \cdots, k_{s}$ are the corresponding degrees of homogeneity, then $\Sigma k_{i}=r$ and, when convenient, we may arrange the notation so that $k_{1} \geq k_{2} \geq \cdots \geq k_{s}$.

If we now replace $e_{1}$ in $F$ by $e_{1}+\lambda e_{i}(i>s)$, the coefficient $H$ of $\lambda$ is not 0 , since $i>s$, and $H$ becomes $k_{1} F$ when $e_{1}$ is replaced by $e_{1}$; it therefore forms a generator of $\mathfrak{F}$ in which the degree of $e_{1}$ is one less than before. It follows that, when $r \leq n$, we may choose a generator which is linear and homogeneous in $r$ units $e_{1}, e_{2}, \cdots, e_{r}$. It is also readily shown that such a tensor defines an irreducible transformable linear set if, and only if, it forms an irreducible set when the transformations of the units are restricted to permuting the first $r e$ 's among themselves. Further, since the choice of fundamental units is arbitrary, we may replace them by variable vectors $x_{1}, x_{2}, \cdots, x_{r}$. For instance, the transformable sets associated with $\Pi_{r}, P_{r}$ and $C_{r}$ are $x_{1} x_{2} \cdots x_{r}$, $\left\{x_{1} x_{2} \cdots x_{r}\right\}$ and $\left|x_{1} x_{2} \cdots x_{r}\right|$, respectively, and of these the first is reducible and the other two irreducible.
5.19 It is not difficult to calculate directly the irreducible transformable sets for small values of $r$ by the aid of the results of the preceding paragraph. If we denote $x_{1}, x_{2}, \cdots$ by $1,2, \cdots$, the following are generators for $r=2,3$.

|  | generator | $r=2$ | order <br> 2.1 |
| :--- | :--- | :---: | :---: |
| $\{12\}$ |  | $n(n+1) / 2$ |  |
| 2.2 | $\|12\|$ |  | $n(n-1) / 2$ |
|  |  | $r=3$ |  |
| 3.1 | $\{123\}$ |  | $n(n+1)(n+2) / 6$ |
| 3.2 | $\|1\| 23\|\mid$ |  | $n\left(n^{2}-1\right) / 3$ |
| 3.3 | $\|1\{23\}\|$ |  | $n\left(n^{2}-1\right) / 3$ |
| 3.4 | $\|123\|$ |  | $n(n-1)(n-2) / 6$. |

This method of determining the generators directly is tedious and the following method is preferable. ${ }^{7}$ Any generator has the form

$$
w_{1}=\Sigma \omega_{i_{1} i_{2}} \cdots i_{r} x_{i_{1}} x_{i_{2}} \cdots x_{i_{r}}
$$

and if $q_{i_{1}} \ldots i_{r}$ denotes the substitution $\left(\begin{array}{llll}1, & 2, & \cdots & r \\ i_{1} & i_{2} & \cdots & i_{r}\end{array}\right)$, we may write

$$
\begin{aligned}
w_{1} & =\Sigma \omega_{i_{1} i_{2}} \ldots i_{i_{r}} q_{i_{1} i_{2}} \ldots i_{r} x_{1} x_{2} \cdots x_{r} \\
& =q_{1}\left(x_{1} x_{2} \cdots x_{r}\right)
\end{aligned}
$$

where $q_{1}$ may be regarded (see chap. 10) as an element of the algebra $S$ whose units are the operators $q$ of the symmetric group on $r$ letters. Now $w_{1}$ generates a transformable set and hence, if $w_{i}=q_{i}\left(x_{1} \cdots x_{r}\right)(i=1,2, \cdots)$ is a
${ }^{7}$ Fuller details of the actual determination of the generators will be found in Weyl: -Gruppentheorie und Quantentheorie, 2 ed. chap. 5.
basis of the set, and $Q$ is the set of elements $q_{1}, q_{2}, \cdots$ in $S$, then the set of elements $Q q=\left(q_{1} q, q_{2} q, \cdots\right)$ must be the same as the set $Q$, that is, in the terminology of chapter $10, Q$ is a semi-invariant subalgebra of $S$; conversely any such semi-invariant subalgebra gives rise to a transformable set and this set is irreducible if the semi-invariant subalgebra is minimal, that is, is contained in no other such subalgebra.

It follows now from the form derived for a group algebra such as $S$ that we get all independent generators as follows. In the first place the operators of $S$ can be divided into sets ${ }^{8} S_{k}(k=1,2, \cdots, t)$ such that (i) the product of an element of $S_{k}$ into an element of $S_{i}(k \neq j)$ is zero; (ii) in the field of complex numbers a basis for each $S_{k}$ can be chosen which gives the algebra of matrices of order $n_{k}^{2}$; and in an arbitrary field $S$ is the direct product of a matric algebra and a division algebra; (iii) there exists a set of elements $u_{k 1}, u_{k 2}, \cdots, u_{k r_{k}}$ in $S_{k}$ such that $\sum_{i} u_{k i}$ is the identity of $S_{k}$ and $u_{k i}^{2}=u_{k i} \neq 0, u_{k i} u_{k j}=0(i \neq j)$ and such that the set of elements $u_{k i} S_{k} u_{k i}$ is a division algebra, which in the case of the complex field contains only one independent element; (iv) the elements of $S_{k}$ can be divided into $\nu_{k}$ sets $u_{k i} S_{k}(i=1,2, \cdots)$ each of which is a minimal semi-invariant subalgebra of $S$ and therefore corresponds to an irreducible transformable set.

[^1]
[^0]:    ${ }^{2}$ The term 'hypernumber' is used in place of vector, as defined in $\$ 1.01$ since we now wish to use the term 'vector' in a more restricted sense.
    ${ }^{2}$ This product was called by Grassmann the general or indeterminate product.

[^1]:    ${ }^{8}$ It is shown in the theory of groups that $t$ equals the number of partitions of $r$.

