# GROUPS OF LINEAR OPERATORS DEFINED BY GROUP CHARACTERS 

## BY

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#### Abstract

Some of the recent work on invariance questions can be regarded as follows: Characterize those linear operators on Hom ( $V, V$ ) which preserve the character of a given representation of the full linear group. In this paper, for certain rational characters, necessary and sufficient conditions are described that ensure that the set of all such operators forms a group $\mathcal{P}$. The structure of $\mathcal{P}$ is also determined. The proofs depend on recent results concerning derivations on symmetry classes of tensors.


1. Statements. Let $G$ be any subgroup of the full linear group $G L(n, C)$ over the complex numbers, and let $\because$ denote the linear closure of $G$ in the total matrix algebra $M_{n}(\mathbf{C})$. Let $K: G \rightarrow \mathrm{GL}(N, \mathbf{C})$ be a representation which is extended to a representation of the multiplicative semigroup of $\mathcal{U}$ in $M_{n}(C)$. Let $\mu_{K}(X)=\operatorname{tr} K(X)$ be the corresponding character. Next, let $\mathscr{L}(G, K)$ denote the multiplicative semigroup of all linear transformations $\mathfrak{T}: \mathfrak{Z} \rightarrow \mathfrak{Q}$ having the property that $\mathfrak{J}$ preserves the character of the representation $K$; that is,

$$
\begin{equation*}
\mu_{K}(\mathscr{T}(X))=\mu_{K}(X), \quad X \in \mathfrak{V} . \tag{I}
\end{equation*}
$$

The two central questions which will concern us in this paper are:
(i) Under what circumstances is $\mathcal{L}(G, K)$ a group, i.e., under what circumstances is it true that if (1) holds, then $\mathfrak{T}$ is nonsingular?
(ii) If $\mathscr{L}(G, K)$ is a group, then what is its structure?

Probably the first instance of a question of this kind was discussed by Frobenius [3] who proved that if $G=G L(n, C)$, so that $\mathscr{Q}=M_{n}(C)$, and if $K(X)=$ $\operatorname{det}(X)$, then $\mathscr{L}(G, K)$ is a group. He proved, in answer to question (ii), that for $\mathfrak{T} \in \mathscr{L}(G, K)$ there exist fixed matrices $U$ and $V$ in $\operatorname{GL}(n, C)$ such that

$$
\mathscr{T}(X)=U X V, \quad X \in M_{n}(\mathrm{C}), \quad \text { or } \quad \mathscr{T}(X)=U X^{T} V, \quad X \in M_{n}(\mathbf{C})
$$

[^0]where $\operatorname{det}(U V)=1\left(X^{T}\right.$ is the transpose of $\left.X\right)$.
A related problem was discussed by I. Schur [12]. Let $3 \leq m \leq n$ and $\mathfrak{T}$ : $M_{n}(\mathrm{C}) \rightarrow M_{n}(\mathrm{C})$ be a linear transformation satisfying the following condition. For each $X \in M_{n}(\mathrm{C})$, the $m$ th order subdeterminants of $\mathscr{T}(X)$ are fixed linearly independent linear homogeneous functions of the $m$ th order subdeterminants of $X$. Schur proved that for such a $\mathfrak{T}$ there exist fixed matrices $U, V \in \mathrm{GL}(n, \mathrm{C})$ such that
$$
\mathscr{T}(X)=U X V, \quad X \in M_{n}(\mathbf{C}), \quad \text { or } \quad \mathscr{T}(X)=U X^{\top} V, \quad X \in M_{n}(\mathbf{C})
$$

This problem can be reformulated in terms of the $m$ th Grassmann compound $C_{m}(X)$ of $X$. Let $S$ be a nonsingular linear transformation from $M_{(\underset{m}{n})}^{(\mathrm{C}) \text { to itself. Char- }}$ acterize those linear transformations $\mathfrak{T}$ on $M_{n}(\mathbb{C})$ which satisfy

$$
C_{m}(\mathscr{T}(X))=S\left(C_{m}(X)\right), \quad X \in M_{n}(\mathrm{C})
$$

This reformulation and a proof depending on more recent results appear in [9].
Let $P=\left[\begin{array}{cc}0 & 1 \\ 1 & 0\end{array}\right]$, and let $G$ be the group consisting of all $X \in \mathrm{GL}(2, \mathrm{C})$ for which $X^{*}=P X^{T} P$. Let $K(X)=\operatorname{det}(X)$. As can be readily verified, $\mathscr{L}(G, K)$ is isomorphic to the set of linear transformations mapping the real space $\mathbf{R}^{4}$ into itself and holding fixed the quadratic form

$$
f(x)=x_{1}^{2}+x_{2}^{2}-x_{3}^{2}-x_{4}^{2}
$$

In $[8]$ it is proved that 9 consists of all $X \in M_{2}(C)$ of the form

$$
\left[\begin{array}{ll}
z & w \\
\bar{w} & z
\end{array}\right]
$$

and that $\mathscr{L}(G, K)$ consists of all $\mathfrak{T}$ of the form

$$
\mathscr{T}(X)=U X V, \quad X \in M_{n}(\mathbf{C}), \quad \text { or } \quad \mathscr{T}(X)=U X^{T} V, \quad X \in M_{n}(\mathbf{C})
$$

where $\operatorname{det}(U V)=1$, and $U^{*}=P U^{T} P, V^{*}=P V^{T} P$.
In [10] it is proved that $\mathscr{L}(G, K)$ is a group when $G=G L(n, C)$ and $K(X)=$ $C_{m}(X)$ for $3<m \leq n$. In this instance, $\mu_{K}(X)=\operatorname{tr} C_{m}(X)$ is the $m$ th elementary symmetric function of the eigenvalues of $X$, or equivalently, the sum of all $\binom{n}{m} m$-square principal subdeterminants of $X$. In case $m<n$, the group $\mathscr{L}\left(\mathrm{GL}(n, \mathbf{C}), C_{m}\right)$ consists of precisely those linear transformations $\mathcal{T}$ of the form

$$
\mathscr{T}(X)=U X V, \quad X \in M_{n}(\mathbf{C}), \quad \text { or } \quad \mathscr{T}(X)=U X^{T} V, \quad X \in M_{n}(\mathbf{C})
$$

where $U V=e^{i \phi} I_{n}$ and $m \phi \equiv O(2 \pi)$. This result was recently extended to include the case $m=3$ in [2].

Another minor modification of our problem occurs in [5]. Let $U(n, \mathrm{C})$ denote the subgroup of $\mathrm{GL}(n, \mathrm{C})$ consisting of all unitary matrices. Then the semigroup
$\mathfrak{£}$ of all linear transformations $\mathscr{T}$ on $M_{n}(\mathrm{C})$ which satisfy $\mathscr{T}(U(n, \mathrm{C})) \subset U(n, \mathrm{C})$ is a group. It is shown that $\mathfrak{J} \in \mathscr{L}$ if and only if there exist fixed matrices $U, V \in$ $U(n, C)$ such that

$$
\mathscr{T}(X)=U X V, \quad X \in M_{n}(\mathrm{C}), \quad \text { or } \quad \mathscr{T}(X)=U X^{T} V, \quad X \in M_{n}(\mathbf{C})
$$

Observe that $\mathscr{L}(\{, K)$ is not always a group. Let $G=G L(n, \mathrm{C})$ and $K(X)$ be the $m$ th Kronecker power $\Pi^{m}(X)$ of $X$ [14]. Then $\mu_{K}(X)=(\operatorname{tr}(X))^{m}$. The annihilator map $\mathfrak{T}$ which sends each $X=\left(x_{i j}\right)$ into $\mathscr{T}(X)=\left(y_{i j}\right)$ where $y_{i j}=\delta_{i j} x_{i j}$ and which clearly belongs to $\mathscr{L}(\not(, K)$, has no inverse.

In this paper we shall discuss problems (i) and (ii) for a certain class of rational representations of the multiplicative semigroup of $\geqslant$ which are in fact components of the mth Kronecker product representation $\Pi^{m}(X)$. It is somewhat easier to state our results in an invariant setting.

Let $H$ be a subgroup of the symmetric group of degree $m$ and let $\chi$ be a character of $H$ of degree 1. Let $V$ be an $n$-dimensional vector space over $\mathbf{C}$; let $U$ be any vector space over $\mathbf{C}$ and $\phi\left(v_{1}, \cdots, v_{m}\right)$ an $m$-multilinear function on the Cartesian product $X_{1}^{m} V$ to $U$. Then $\phi$ is said to be symmetric with respect to $H$ and $\chi$ if

$$
\phi\left(v_{\sigma(1)}, \cdots, v_{\sigma(m)}\right)=\chi(\sigma) \phi\left(v_{1}, \cdots, v_{m}\right)
$$

holds for any $\sigma \in H$ and arbitrary vectors $v_{i} \in V$. A pair ( $P, \nu$ ) consisting of a vector space $P$ over $C$ and a fixed $m$-multilinear function $\nu: X_{1}^{m} V \rightarrow P$, symmetric with respect to $H$ and $\chi$, is a symmetry class of tensors associated with $H$ and $\chi$ if
(i) $\langle\boldsymbol{\operatorname { n g }} \nu\rangle=P$, i.e., the linear closure of the range of $\nu$ is $P$;
(ii) (universal factorization property) for any vector space $U$ over $\mathbf{C}$ and any $m$-multilinear function $\phi: X_{1}^{m} V \rightarrow U$, symmetric with respect to $H$ and $\chi$, there exists a unique linear function $b: P \rightarrow U$ such that $\phi=b \nu$; i.e., the following diagram is commutative.


For any linear transformation $X: V \rightarrow V$ the preceding universal factorization property permits us to define a unique linear transformation $K(X): P \rightarrow P$, the induced transformation on $P$, which satisfies the following identity. For arbitrary vectors $v_{1}, \cdots, v_{m}$ in $V$

$$
\begin{equation*}
K(X) \nu\left(v_{1}, \cdots, v_{m}\right)=\nu\left(X v_{1}, \cdots, X v_{m}\right) . \tag{2}
\end{equation*}
$$

By the spanning property of the range of $\nu$ (i.e., (i) above), (2) immediately implies that $K(X)$ is multiplicative and in fact if $m \leq n$ and $X \in \mathrm{GL}_{n}(V)$, the group of all linear bijections on $V$, then $K(X) \in \mathrm{GL}_{N}(P)$ where $N=\operatorname{dim} P$. If $G$ is any subgroup of $\mathrm{GL}_{n}(V)$ and $थ$ is the linear closure of $G$ in $\operatorname{Hom}(V, V)$, we are thus in a position to discuss the structure of $\mathcal{L}(G, K)$, which for the class of representations $K(X)$ just defined depends on the group $H$ and the character $\chi$. If we identify $V$ with the space of $n$-tuples over C , then of course $\mathrm{GL}_{n}(V)$ can be identified with GL ( $n, \mathrm{C}$ ) and we can ask for the structure of the semigroup $\mathscr{L}(G, K)$ for the preceding class of representations $K(X)$ of $थ$.

Our main results follow.
Theorem 1. Let $\operatorname{dim} V=n, H \subset S_{m}, \chi$ a character of degree 1 on $H$. Let $(P, \nu)$ be the symmetry class associated with $H$ and $\chi$ and $X \rightarrow K(X)$ be a representation of $G=G L_{n}(V)$ by induced transformations on $P$. If $m \leq n$ or $\chi \equiv 1$, then $\mathscr{L}\left(\mathrm{GL}_{n}(V), K\right)$ is a group if and only if $H \neq\{e\}$.

Theorem 2. Let $\operatorname{dim} V=n, H=S_{m^{\prime}} m>1, \chi \equiv 1$. Let $G$ be a subgroup of $\mathrm{GL}_{n}(V)$. If the algebra $\mathfrak{N}$ bas the property that the conjugate transpose $X^{*}$ of each $X$ in I is again in 2 , then $\mathscr{L}(G, K)$ is a group.

Theorem 3. In Theorem 1, take $H=S_{m}, m \geq 3$ and $\chi \equiv 1$. Let $\mathscr{L}_{1}\left(\mathrm{GL}_{n}(V), K\right)$ denote the subgroup of $\mathscr{L}\left(\mathrm{GL}_{n}(V), K\right)$ of those $\mathfrak{T}: \operatorname{Hom}(V, V) \rightarrow \operatorname{Hom}(V, V)$ satisfying $\mathscr{T}\left(I_{V}\right)=\xi I_{V}$. Then $\mathscr{\bigotimes}_{1}^{n}\left(\mathrm{GL}_{n}(V), K\right)$ consists precisely of those linear transformations $\mathfrak{T}$ which have the form

$$
\begin{equation*}
\mathscr{T}(X)=\xi U^{-1} X U, \quad X \in \operatorname{Hom}(V, V), \tag{3}
\end{equation*}
$$

or

$$
\begin{equation*}
\mathscr{T}(X)=\xi U^{-1} X^{T} U, \quad X \in \operatorname{Hom}(V, V) \tag{4}
\end{equation*}
$$

Theorem 4. In Theorem 1, take $H=A_{m} \subset S_{m}$ to be the alternating group, $m \geq 3$, and $\chi \equiv 1$. The group $£_{1}\left(\mathrm{GL}_{n}(V), K\right)$ consists precisely of those linear transformations $\mathfrak{J}$ of the form (3) or (4).

Corollary 1. Let $\operatorname{dim} V=n, H=S_{m}, m>1, \chi \equiv 1$. If $G$ is the group of all $n \times n$ permutation matrices (so that $\mathcal{U}$ is the algebra of generalized doubly stocbastic matrices), then $\mathscr{L}(G, K)$ is a group.

Let $m$ and $n$ be positive integers. Let $Q_{m, n}$ (resp. $G_{m, n}$ ) denote the set of all strictly increasing (resp. nondecreasing) sequences of length $m$ chosen from the set $\{1,2, \cdots, n\}$. If $f\left(\lambda_{1}, \cdots, \lambda_{n}\right)$ is a polynomial symmetric in the indeter-
minates $\lambda_{1}, \cdots, \lambda_{n}$ and $X \in \operatorname{Hom}(V, V)$, we shall denote by $f(X)$ the value of $f$ at the eigenvalues of $X$. For $m \geq 1$, let $b_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)$ denote the $m$ th completely symmetric polynomial

$$
b_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)=\sum_{\alpha \in G_{m, n}} \prod_{t=1}^{m} \lambda_{\alpha(t)}
$$

and let $k_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)$ denote the symmetric polynomial

$$
\begin{aligned}
k_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right) & =b_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)+\sum_{a \in Q_{m, n}} \prod_{t=1}^{m} \lambda_{\alpha(t)} \\
& =b_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)+E_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)
\end{aligned}
$$

where $E_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)$ is the $m$ th elementary symmetric function of $\lambda_{1}, \cdots, \lambda_{n}$ when $m \leq n$ and 0 if $m>n$.

Corollary 2. Let $m \geq$ 3. Any linear transformation $\mathfrak{T}: \operatorname{Hom}(V, V) \rightarrow$ $\operatorname{Hom}(V, V)$ satisfying $\mathscr{J}\left(I_{V}\right)=\xi I_{V}$ and $b_{m}(\mathscr{T}(X))=b_{m}(X), X \in \operatorname{Hom}(V, V)$, bas the form (3) or (4).

Corollary 3. Let $m \geq$ 3. Any linear transformation $\mathfrak{T}: \operatorname{Hom}(V, V) \rightarrow$ $\operatorname{Hom}(V, V)$ satisfying $\mathscr{T}\left(I_{V}\right)=\xi I_{V}$ and $k_{m}(T(X))=k_{m}(X), X \in \operatorname{Hom}(V, V)$, bas the form (3) or (4).

We conjecture that in fact $\mathscr{L}_{1}\left(\mathrm{GL}_{n}(V), K\right)=\mathscr{L}\left(\mathrm{GL}_{n}(V), K\right)$ in Theorems 3 and 4. This amounts to showing that if $\mathfrak{T}: \mathrm{GL}_{n}(V) \rightarrow \mathrm{GL}_{n}(V)$ and $\mu_{K}(\mathscr{T}(X))=\mu_{K}(X)$ holds for all $X \in \operatorname{Hom}(V, V)$, then $\mathscr{T}\left(I_{V}\right)=\xi I_{V}$ where $\xi^{m}=1$.
2. Partial derivations. In [6] the standard notion of a derivation on a tensor algebra [4] is extended to higher order derivations on a general symmetry class ( $P, \nu$ ). We shall further extend the idea of a derivation induced by a single linear transformation to partial derivations induced by two linear transformations [11].

Let $T$ and $S$ be in Hom $(V, V)$ and let $r+s=m$. For $\omega \in Q_{r, m}$ define

$$
\begin{equation*}
\Pi_{\omega}(T, s)=\bigotimes_{i=1}^{m} X_{i} \tag{5}
\end{equation*}
$$

where $X_{i}=T$ for $i \in \mathrm{mg} \omega$ and $X_{i}=S$ otherwise. In other words (5) is the tensor product of the linear transformations $T$ and $S$ in which $T$ appears in positions numbered $\omega$ and $S$ appears elsewhere. The linear transformation (5) acts on $\bigotimes_{1}^{m} V$, which of course is the symmetry class associated with $H=\{e\}$. Define

$$
\delta_{r, s}(T, S)=\sum_{\omega \in Q_{r, m}} \Pi_{\omega}(T, S)
$$

In order to simplify subsequent notation we make the following convention. Let $f: Q_{r, m} \times Q_{s, m} \rightarrow R$ be any function into a set $R$ having an associative addition. We shall let $\Sigma^{\prime} f(\omega)$ denote the summation of $f(\omega, \gamma)$ over all sequences $\omega \epsilon$ $Q_{r, m^{\prime}} \gamma \in Q_{s, m}$ such that mg $\omega \cap \mathrm{mg} \gamma=\varnothing$. Next, define

$$
M\left(X_{1}, \cdots, X_{m}\right)=\sum_{\phi \in S_{m}} X_{\dot{\phi}(1)} \otimes \cdots \otimes X_{\phi(m)} .
$$

Then it is easy to show that

$$
\begin{equation*}
M_{r, s}(T, S)=r!s!\delta_{r, s}(T, S) \tag{6}
\end{equation*}
$$

where $M_{r, s}(T, S)$ denotes

$$
\overbrace{M(T, \cdots, T}^{r}, \frac{s}{S, \cdots, S}) .
$$

It is also a standard fact concerning symmetry classes that if the symmetry operator associated with $H$ and $\chi$ (a linear transformation on $\otimes_{1}^{m} V$ ) is defined by

$$
\tau_{x}=\frac{1}{|H|} \sum_{\sigma \in H} \chi(\sigma) \sigma
$$

$\left(\sigma\left(v_{1} \otimes \ldots \otimes v_{m}\right)=v_{\sigma-1(1)} \otimes \ldots \otimes_{v_{-1} 1_{(m)}}\right)$, then the pair $(P, \nu), P=\operatorname{rng} \tau_{\chi}$ $C \otimes_{1}^{m} V, \nu\left(v_{1}, \cdots, v_{m}\right)=r_{x}\left(\nu_{1} \otimes \ldots \otimes v_{m}\right)$ is the symmetry class associated with $H$ and $\chi$. It is also easy to show that the transformation $M_{r, s}(T, S)$ satisfies $M_{r, s}(T, S)_{\sigma}=\sigma M_{r, s}(T, S)$ for all $\sigma \in S_{m}$; hence any symmetry class is an invariant subspace of $M_{r, s}(T, S)$. But then in view of (6), each symmetry class is an invariant subspace of $\delta_{r, s}(T, S)$.

We define the $(r, s)$ partial derivation associated with $T$ and $S$ on ( $P, \nu$ ) to be the restriction of $\delta_{r, s}(T, S)$ to the invariant subspace $P$. We denote this by $\Omega_{r, s}(T, S)$. The reason for calling $\Omega_{r, s}(T, S)$ the ( $\left.r, s\right)$ partial derivation on ( $P, \nu$ ) is the following formula:

$$
\begin{equation*}
K\left(x_{1} T+x_{2} S\right)=\sum_{r+s=m} x_{1}^{r} x_{2}^{s} \Omega_{r, s}(T, S) \tag{7}
\end{equation*}
$$

In order to verify (7) we compute that
$K\left(x_{1} T+x_{2} S\right) \nu\left(\nu_{1}, \cdots, v_{m}\right)=\nu\left(\left(x_{1} T+x_{2} S\right) v_{1}, \ldots,\left(x_{1} T+x_{2} S\right) \nu_{m}\right)$

$$
\begin{equation*}
=\sum_{r+s=m} x_{1}^{r} x_{2}^{s} \sum^{\prime} \nu\left(\cdots, T v_{\omega(1)}, \cdots, S v_{\gamma(1)}, \cdots, T v_{\omega(r)}, \cdots, S v_{\gamma(s)}, \cdots\right) \tag{8}
\end{equation*}
$$

where in the inside summand on the right side of (8) the $T$ occurs in precisely the positions numbered $\omega$ and the $S$ in positions numbered $\gamma$. On the other hand,
$\Omega_{r, s}(T, S) \mathcal{\nu}\left(v_{1}, \cdots, v_{m}\right)=\delta_{r, s}(T, S)_{\tau_{\chi}}\left(v_{1} \otimes \cdots \otimes v_{m}\right)=\tau_{\chi} \delta_{r, s}(T, S)\left(v_{1} \otimes \cdots \otimes v_{m}\right)$

$$
\begin{align*}
& =r_{X} \sum^{\prime}\left(\cdots \otimes T v_{\omega(1)} \otimes \cdots \otimes S v_{\gamma(1)} \otimes \cdots \otimes T v_{\omega(r)} \otimes \cdots \otimes S v_{\gamma(s)} \otimes \cdots\right)  \tag{9}\\
& =\sum^{\prime} \nu\left(\cdots, T v_{\omega(1)}, \cdots, S v_{\gamma(1)}, \cdots, T v_{\omega(r)}, \cdots, S v_{\gamma(s)}, \cdots\right)
\end{align*}
$$

Replacing (9) in (8) we have (7).
We observe a number of elementary facts concerning the partial derivation $\Omega_{r, s}(T, S):$
(i) If $X$ and $Y$ are in $\operatorname{Hom}(V, V)$, then

$$
\begin{equation*}
K(X) \Omega_{r, s}(T, S) K(Y)=\Omega_{r, s}(X T Y, X S Y) \tag{10}
\end{equation*}
$$

This follows immediately from (7).
(ii) If $V$ is a unitary space, then $\bigotimes_{1}^{m} V$ is also a unitary space. Thus there is a natural inner product induced on the symmetry class ( $P, \nu$ ) associated with $H$ and $\chi$. Moreover, if $T^{*}$ is the conjugate dual of $T \in \operatorname{Hom}(V, V)$, then the conjugate dual of $\Omega_{r, s}(T, S)$ with respect to the induced inner product in ( $P, \nu$ ) is

$$
\begin{equation*}
\Omega_{r, s}(T, S)^{*}=\Omega_{r, s}\left(T^{*}, S^{*}\right) \tag{11}
\end{equation*}
$$

(iii) $\Omega_{1, m-1}(T, S)$ is linear in $T$ and $\Omega_{m-1,1}(T, S)$ is linear in $S$.

A somewhat more combinatorially involved description is necessary to itemize the eigenvalues of $\Omega_{r, s}(T, S)$. In order to describe a basis for an arbitrary symmetry class assaciated with $H$ and $\chi$, we regard the elements of $H$ as permutations acting on the functions (i.e., sequences) in $\Gamma_{m, n}=Z_{n}^{Z_{m}}$, where $Z_{q}=$ $\{1,2, \cdots, q\}$ and for $\sigma \in H, a \in \Gamma_{m, n}$,

$$
\sigma((\alpha))(t)=\alpha\left(\sigma^{-1}(t)\right), \quad t \in Z_{m}
$$

Let $\Delta$ denote a system of distinct representatives for the orbits in $\Gamma_{m, n}$ induced by $H$, and let $\bar{\Delta}$ denote the set of all of those elements $\alpha \in \Delta$ for which the character $\chi$ is identically 1 on the stabilizer subgroup $H_{\alpha}=\{\sigma \in H: \sigma(\alpha)=\alpha\}$. Let $n(\alpha)=\left|H_{a}\right|$. It is routine to verify that if $\left\{e_{1}, \ldots, e_{n}\right\}$ is a basis of $V$, then the decomposable elements $\nu\left(e_{\alpha(1)}, \cdots, e_{a(m)}\right), \alpha \in \bar{\Delta}$ form a basis for $P$. In fact, if $\left\{e_{1}, \cdots, e_{n}\right\}$ is an orthonormal (hereafter abbreviated o.n.) basis of $V$, then the $|\bar{\Delta}|$ decomposable elements $(|H| / n(\alpha))^{1 / 2} \nu\left(e_{\alpha(1)}, \cdots, e_{\alpha(m)}\right)$ form an o.n. basis for $P$ with respect to the induced inner product in $\bigotimes_{1}^{m} V$ defined by

$$
\left(x_{1} \otimes \cdots \otimes x_{m}, y_{1} \otimes \cdots \otimes y_{m}\right)=\prod_{i=1}^{m}\left(x_{i}, y_{i}\right)
$$

If we choose the system of distinct representatives $\Delta$ so that each sequence $\alpha \in \Delta$ is lowest in lexicographic order in the orbit in which it lies, then it is easy to see that $Q_{m, n} \subset \bar{\Delta}$ and $G_{m, n} \subset \Delta$, whatever the group $H$ and character $\chi$ may be [7].

If we use the fact that for any pair ( $T, S$ ) of commuting linear transformations there exists a common triangular o.n. basis, then it is not difficult to prove the. following [11]:
(iv) If $S T=T S$ and the eigenvalues of $T$ and $S$ are $\lambda_{1}, \cdots, \lambda_{n}$ and $\kappa_{1}, \cdots, \kappa_{n}$ respectively, then after a suitable reordering of the $\kappa_{i}$ 's, the eigenvalues of $\Omega_{r, s}(T, S)$ are the numbers

$$
\begin{equation*}
\sum^{\prime} \prod_{i=1}^{\tau} \lambda_{\alpha \omega(i)} \prod_{j=1}^{s} \kappa_{\alpha \gamma(j)} \quad a \in \bar{\Delta} \tag{12}
\end{equation*}
$$

In particular, if $r=1$ and $s=m-1$, then the eigenvalues of $\Omega_{1, m-1}(T, S)$ are

$$
\begin{equation*}
\sum_{t=1}^{m} \lambda_{\alpha(t)} \prod_{j \sim t}^{m} \kappa_{\alpha(j)}, \quad \alpha \in \bar{\Delta} \tag{13}
\end{equation*}
$$

Some additional combinatorial maneuvering will be required: if $\alpha \in \Gamma_{m, n}$ and $1 \leq t \leq n$, then let $m_{t}(\alpha)$ denote the number of integers $i$ in $\{1,2, \ldots, m\}$ for which $\alpha(i)=t$; i.e., $m_{t}(\alpha)$ is the multiplicity of occurrence of $t$ in the range of $\alpha$. More generally, if $p_{1}+\cdots+p_{r}=n$ is a partition of $n$ into positive parts, we define

$$
\eta_{t}(\alpha)=\sum_{j=P_{t-1}+1}^{P_{t}} m_{j}(\alpha)
$$

where $P_{t}=p_{1}+\cdots+p_{t}$; i.e., $\eta_{t}(\alpha)$ is the number of times any integer $k$ satisfying $P_{t-1}<k \leq P_{t}$ occurs in the range of $a$. We can write the eigenvalues (13) of $\Omega_{1, m-1}(T, S)$ in a form somewhat more suitable for our subsequent computations. Suppose that the eigenvalues of $T$ are given by

$$
\begin{equation*}
\lambda_{1}=\cdots=\lambda_{P_{1}}=l_{1} ; \lambda_{P_{1}+1}=\cdots=\lambda_{P_{2}}=l_{2} ; \cdots ; \lambda_{P_{r-1}+1}=\cdots=\lambda_{P_{r}}=l_{r} \tag{14}
\end{equation*}
$$

where the numbers $l_{t}$ are distinct. Suppose, moreover, that $f(x)$ is an arbitrary scalar polynomial and $S=f(T)$. In this case the ordering (14) of the eigenvalues of $T$ induces a corresponding ordering of the eigenvalues $\kappa_{1}, \cdots, \kappa_{n}$ of $S$, i.e.,

$$
\kappa_{1}=\cdots=\kappa_{P_{1}}=k_{1} ; \kappa_{P_{1}+1}=\cdots=\kappa_{P_{2}}=k_{2} ; \cdots ; \kappa_{P_{r-1}+1}=\cdots=\kappa_{n}=k_{r}
$$

Since the ploynomial $f(x)$ can be chosen arbitrarily, it follows that the numbers $k_{t}$ may be chosen arbitrarily. Regarding the $k_{t}$ as momentarily all different from zero, we see that the eigenvalues (13) become

$$
\begin{align*}
\sum_{t=1}^{m} \lambda_{a(t)} \prod_{j \neq t}^{m} \kappa_{\alpha(j)} & =\sum_{t=1}^{m} \frac{\lambda_{a(t)}}{\kappa_{a(t)}} \prod_{j=1}^{m} \kappa_{a(j)}=\sum_{t=1}^{m} \frac{\lambda_{a(t)}}{\kappa_{a(t)}} \prod_{j=1}^{n} \kappa_{j}^{m_{j}(a)} \\
& =\sum_{t=1}^{n} m_{t}(\alpha) \frac{\lambda_{t}}{\kappa_{t}} \prod_{j=1}^{n} \kappa_{j}^{m_{j}(\alpha)}=\sum_{t=1}^{r} \eta_{t}(a) \frac{l_{t}}{k_{t}} \prod_{j=1}^{r} k_{j}^{\eta_{j}(a)}  \tag{15}\\
& =\sum_{t=1}^{r} \eta_{t}(a) l_{t} k_{t}^{\eta_{t}(a)-1} \cdot \prod_{j \neq t}^{r} k_{j}^{\eta_{j}^{(a)}}
\end{align*}
$$

If we interpret $0^{0}$ as 1 , then (15) holds even when some of the numbers $k_{t}$ are zero. We now have the following formula for the trace of $\Omega_{1, m-1}(T, S)$ :

$$
\begin{align*}
\operatorname{tr} \Omega_{1, m-1}(T, S) & =\sum_{a \in \bar{\Delta}} \sum_{t=1}^{r} \eta_{t}(\alpha) l_{t} k_{t}^{\eta_{t}(\alpha)-1} \prod_{j \neq t}^{r} k_{j}^{\eta_{j}^{(\alpha)}} \\
& =\sum_{t=1}^{r} l_{t}\left(\sum_{a \in \bar{\Delta}} \eta_{t}(\alpha) k_{t}^{\eta_{t}(\alpha)-1} \prod_{j \neq t}^{r} k_{j}^{\eta_{j}^{(\alpha)}}\right) \tag{16}
\end{align*}
$$

## 3. Proofs.

Lemma 1. Let $G$ be any subgroup of $\mathrm{GL}_{n}(V)$ and $\mathcal{T} \in \mathcal{L}(G, K)$. Then if $A$ $\epsilon$ ker $\boldsymbol{J}$,

$$
\begin{equation*}
\operatorname{tr} \Omega_{1, m-1}(A, X)=0, \quad X \in \mathfrak{U} . \tag{17}
\end{equation*}
$$

Proof. Let $x_{1}, x_{2}$ be indeterminates over C. From (7) we have

$$
K\left(x_{1} A+x_{2} X\right)=\sum_{r=0}^{m} x_{1}^{r} x_{2}^{m-r} \Omega_{r, m-r}(A, X) .
$$

Thus, if $\mathscr{T}(A)=0$ we have

$$
\begin{align*}
\operatorname{tr} \sum_{r=0}^{m} x_{1}^{r} x_{2}^{m-r} \Omega_{r, m-r}(A, X) & =\mu_{K}\left(x_{1} A+x_{2} X\right)=\mu_{K}\left(\mathscr{T}\left(x_{1} A+x_{2} X\right)\right) \\
& =\mu_{K}\left(\mathscr{T}\left(x_{2} X\right)\right)=\mu_{K}\left(x_{2} X\right)=x_{2}^{m} \mu_{K}(X) . \tag{18}
\end{align*}
$$

If we equate coefficients in (18) we obtain

$$
\operatorname{tr} \Omega_{r, m-r}(A, X)=0, \quad r=1,2, \cdots, m,
$$

and hence (17) follows.
Proof of Theorem 1. Assume $H \neq\{e\}$ and let $A \in$ ker $\mathfrak{T}$. We show that $A=$ 0 . By (10) we can assume that (17) holds for all $X$ and that $A$ is in Jordan normal form.

Lemma 2. If every elementary divisor of $A$ is linear and $A \in$ ker $\mathcal{T}$, then $A=0$.

Proof. By (iii) of $\$ 2$, (17) becomes

$$
\begin{equation*}
\sum_{i=1}^{n} a_{i i} \operatorname{tr} \Omega_{1, m-1}\left(E_{i i}, X\right)=0, \quad X \in M_{n}(\mathrm{C}) \tag{19}
\end{equation*}
$$

We first assume that $\chi \equiv 1$. The eigenvalues of $E_{i i}$ are $\lambda_{i}=1$ and $\lambda_{j}=0$, $j \neq i$; while those of $E_{k k}$ are $\kappa_{k}=1$ and $\kappa_{j}=0, j \neq k$. Hence

$$
\begin{equation*}
\operatorname{tr} \Omega_{1, m-1}\left(E_{k k^{\prime}} E_{k k}\right)=\sum_{a \in \Lambda} \sum_{t=1}^{m} \kappa_{a(t)} \prod_{j \neq t}^{m} \kappa_{a(j)}=\sum_{t=1}^{m} \kappa_{k}^{m}=m, \tag{20}
\end{equation*}
$$

because the only term which survives in the inner summation in (20) is the term corresponding to that $a$ for which $\alpha(1)=a(2)=\cdots=a(m)=k$. We remark that this sequence is always in $\bar{\Delta}$ since, as we remarked, $G_{m, n} \subset \bar{\Delta}$ when $\chi \equiv 1$. If $i \neq k$, we again compute that

$$
\begin{equation*}
\operatorname{tr} \Omega_{1, m-1}\left(E_{i i}, E_{k k}\right)=\sum_{\alpha \in \bar{\Lambda}} \sum_{t=1}^{m} \lambda_{\alpha(t)} \prod_{j \neq t} \kappa_{\alpha(j)}=\sum_{\beta \in \bar{\Delta}} 1 \cdot 1^{m-1} \tag{21}
\end{equation*}
$$

where the inner summation in (21) is over precisely those $\beta \in \bar{\Delta}$ for which $m_{k}(\beta)=$ $m-1$ and $m_{i}(\beta)=1$. Once again, since $G_{m, n} \subset \bar{\Delta}$, such sequences exist and we let $p_{i k}$ denote their number. We assert that $p_{i k}$ is independent of the pair ( $i, k$ ); for if $P$ is an arbitrary permutation matrix we have from (10)

$$
p_{i k}=\operatorname{tr} \Omega_{1, m-1}\left(E_{i i^{\prime}}, E_{k k}\right)=\operatorname{tr} \Omega_{1, m-1}\left(P^{T} E_{i i} P, P^{T} E_{k k} P\right)
$$

Obviously if $i \neq k$, we can choose $P$ so that $P^{T} E_{i i} P=E_{i^{\prime}{ }_{i} \prime}$, and $P^{T} E_{k k} P=$ $E_{k^{\prime} k^{\prime}}$ for any preassigned distinct integers $i^{\prime}, k^{\prime}$. We set $p$ equal to the common value of the $p_{i k}$. We next assert that $p<m$, for since $H \neq\{e\}$, there must exist at least $t$ wo sequences $\alpha \neq \beta$ in the same $H$-orbit for which $m_{k}(\alpha)=m_{k}(\beta)=$ $m-1$ and $m_{i}(\alpha)=m_{i}(\beta)=1$. Thus there are at most $m-1$ elements of $\bar{\Delta}$ with that property. If we set $X$ successively equal to $E_{k k^{\prime}} k=1,2, \cdots, n$, in (19) we obtain the following system of linear equations:

$$
m a_{i i}+\sum_{k \neq i} p a_{k k}=0, \quad 1 \leq i \leq n .
$$

Since $p<m$, the coefficient matrix in this system is nonsingular and we conclude that $a_{i i}=0$ for $1 \leq i \leq n$. Thus since the elementary divisors of $A$ are linear we conclude $A=0$.

We now consider the case in which $\chi \not \equiv 1$. If $x$ is an indeterminate, then since $E_{i i}$ and $I_{n}+x E_{11}$ commute we have from (13) that

$$
\begin{align*}
\operatorname{tr} \Omega_{1, m-1}\left(E_{i i}, I_{n}+x E_{11}\right) & =\sum_{a \in \bar{\Delta}} \sum_{t=1}^{m} \lambda_{\alpha(t)} \prod_{j \neq t} \kappa_{\alpha(j)} \\
& =\sum_{a \in \bar{\Phi}} \prod_{j=1}^{n} \kappa_{j}^{m_{j}^{(a)}} \sum_{t=1}^{n} m_{t}(\alpha) \frac{\lambda_{t}}{\kappa_{t}} \tag{22}
\end{align*}
$$

where $\lambda_{1}, \cdots, \lambda_{n}$ are again the eigenvalues of $E_{i i}$ and $\kappa_{1}, \cdots, \kappa_{n}$ are the eigenvalues of $I_{n}+x E_{11}$. In case $i=1$, the right side of (22) becomes

$$
\begin{equation*}
\sum_{a \in \Delta} m_{1}(\alpha)(1+x)^{m_{1}(\alpha)-1} \tag{23}
\end{equation*}
$$

which we denote by $\phi_{1}(x)$. If $i>1$, we see that the value of (22) (which is of course the same for $i=2,3, \cdots, n$ ) is

$$
\begin{equation*}
\phi_{2}(x)=\sum_{a \in \bar{\Delta}} m_{2}(\alpha)(1+x)^{m_{1}(\alpha)} \tag{24}
\end{equation*}
$$

With $X=I_{n}+x E_{11}$ in (19) we have from (23) and (24) that

$$
\begin{align*}
& a_{11} \phi_{1}(x)+\phi_{2}(x) \sum_{i=2}^{n} a_{i i} \\
& =a_{11} \sum_{\alpha \in \Phi} m_{1}(\alpha)(1+x)^{m_{1}(\alpha)-1}+\left(\sum_{i=2}^{n} a_{i i}\right)\left(\sum_{\alpha \in \Phi} m_{2}(\alpha)(1+x)^{m_{1}(\alpha)}\right)=0 . \tag{25}
\end{align*}
$$

We observe that

$$
\begin{equation*}
q_{j}=\sum_{\alpha \in \bar{\Sigma}} m_{j}(\alpha) \tag{26}
\end{equation*}
$$

is independent of $j$ and we denote this common value by $q$. A proof of this can be based on the fact that the eigenvalues of $K(X)$ are

$$
\prod_{j=1}^{m} x_{\alpha(j)}=\prod_{j=1}^{n} x_{j}^{m_{j}^{(\alpha)}}, \quad \alpha \in K
$$

where $x_{1}, \cdots, x_{n}$ are the eigenvalues of $X$. Thus

$$
\operatorname{det}(K(X))=\prod_{\alpha \in \triangle} \prod_{j=1}^{n} x_{j}^{m_{j}^{(\alpha)}}=\prod_{j=1}^{n} x_{j}^{q_{j}}
$$

Since $\operatorname{det}\left(K\left(P^{T} X P\right)\right)=\operatorname{det}(K(X))$ for any permutation matrix $P$, we conclude that

$$
\prod_{j=1}^{n} x_{j}^{q_{j}}=\prod_{j=1}^{n} x_{\sigma(j)}^{q_{j}}, \quad \sigma \in S_{n} .
$$

And since $x_{1}, \cdots, x_{n}$ are arbitrary, $q_{1}=q_{2}=\cdots=q_{n}$. For example, it is easy to compute that

$$
q=q_{1}=\sum_{a \in Q_{m, n}} m_{1}(\alpha)=\binom{n-1}{m-1}
$$

and since $K(X)$ is the familiar $m$ th compound martix $C_{m}(X)$ we see that

$$
\operatorname{det} C_{m}(X)=\prod_{j=1}^{m} x_{j}^{q}=(\operatorname{det}(X))^{\binom{n-1}{m-1}}
$$

and we have as a corollary to our computations the well-known Sylvester-Franke theorem [1]. For $x=0$, (25) becomes $q \operatorname{tr} A=0$ and hence $\operatorname{tr} A=0$; we therefore rewrite (25) as

$$
\begin{equation*}
\left(\phi_{1}(x)-\phi_{2}(x)\right) a_{11}=0 \tag{27}
\end{equation*}
$$

Clearly,

$$
\operatorname{deg} \phi_{1}(x)=\max _{\alpha \in \bar{\Delta}} m_{1}(\alpha)-1 .
$$

Now $\chi \not \equiv 1$, and thus there exists $\beta \in \bar{\Delta}$ and $1<j \leq n$ such that

$$
j \in \operatorname{rng} \beta, \quad m_{1}(\beta)=\max _{a \in \frac{\Delta}{\Delta}} m_{1}(\alpha)
$$

for otherwise $\max _{a \in \bar{\Delta}} m_{1}(\alpha)=m$ and $\alpha=(1,1, \cdots, 1) \in \bar{\Delta}$. But the stabilizer of this $a$ is obviously all of $H$ and it would follow that $\Sigma_{\sigma \epsilon H} \chi(\sigma) \neq 0$ so that $\Sigma_{\sigma \epsilon H} \chi(\sigma)=|H|$. This can only happen if $\chi \equiv 1$, since $\chi$ is a character of degree 1. Hence

$$
\operatorname{deg} \phi_{1}(x)=m_{1}(\beta)-1 \quad \text { and } \quad \operatorname{deg} \phi_{2}(x)=m_{1}(\beta),
$$

so that $\phi_{1}(x)-\phi_{2}(x) \neq 0$. From (27) it follows that $a_{11}$ (and hence any $a_{i i}$ ) is 0 . Thus $A=0$, completing the proof of Lemma 2. We can now remove the condition that $A$ has linear elementary divisors.

Lemma 3. If $A \in \operatorname{ker} \mathscr{T}$, then $A=0$.
Proof. From (11) we know that

$$
\Omega_{1, m-1}(A, X)^{*}=\Omega_{1, m-1}\left(A^{*}, X^{*}\right)
$$

Since (17) holds for all $X$, we have $\operatorname{tr} \Omega_{1, m-1}\left(A^{*}, X\right)=0$ and hence $\operatorname{tr} \Omega_{1, m-1}\left(A \pm A^{*}, X\right)=0$. Now $A \pm A^{*}$ is normal and hence has linear elementary divisors. Applying Lemma 2, we conclude that $A \pm A^{*}=0$ and therefore $A=0$.

We have proved that if $H \neq\{e\}$ and $m \leq n$ or $\chi \equiv 1$, then any $\mathfrak{T} \epsilon$ $\mathscr{L}\left(\mathrm{GL}_{n}(V), K\right)$ is nonsingular; since we have observed that this set obviously forms a semigroup, it is in fact a group. Conversely, if $H=\{e\}$, then the symmetry class $P$ is just the $m$ th tensor space $\bigotimes_{1}^{m} V$ and $K(X)=\Pi^{m}(X)$. As we saw in $\S 1, \mathscr{L}(G, K)$ is not a group. This completes the proof of Theorem 1.

Proof of Theorem 2. We observe that if $A \in \mathbb{Z}$ and $f(x)$ is any scalar polynomial, then $f(A) \in \mathfrak{A}$. The eigenvalues of $\Omega_{1, m-1}(A, f(A))$ are given by (13) where $\lambda_{1}, \cdots, \lambda_{n}$ are the eigenvalues of $A$ and $\kappa_{1}, \cdots, \kappa_{n}$ are the eigenvalues of $f(A)$. Moreover, it is clear that the values of $f(x)$ may be arbitrarily assigned at the distinct eigenvalues of $A$.

Since $H=S_{m}$ and $\chi \equiv 1$, the symmetry class $P$ is precisely the $m$ th completely symmetric space, usually denoted by $V^{(m)}$ [4], and the sequence set $\bar{\Delta}$ is precisely $G_{m, n}$. Once again the problem is to show that if $\mathscr{T}(A)=0$, then $A=0$.

We have from Lemma 1 and formula (16) that

$$
\begin{equation*}
\sum_{t=1}^{\Gamma} l_{t}\left(\sum_{a \in \overline{\mathbf{\Sigma}}} \eta_{t}(\alpha) k_{t}^{\eta_{t}^{(\alpha)-1}} \prod_{j \neq t}^{r} k_{j}^{\eta_{j}^{(\alpha)}}\right)=0 \tag{28}
\end{equation*}
$$

in which the distinct eigenvalues of $A$ are described in (14) and the numbers $k_{1}, \cdots, k_{r}$ may be chosen arbitrarily. The first part of the proof is devoted to showing that all the eigenvalues of $A$ are equal (i.e., that $r=1$ ). Assume then that $r>1$. For a fixed $t$, set $k_{t}=1$ and $k_{j}=0$ for $j \neq t$; observe that the coefficient of $l_{t}$ in (28) is

$$
m \sum_{a \in \bar{\Delta}, \eta_{t}(\alpha)=m} 1=m\binom{p_{t}+m-1}{m}
$$

where the indicated binomial coefficient is just a count of the total number of sequences $a$ in $G_{m, n}$ such that rng $a \subset\left\{P_{t-1}+1, \cdots, P_{t}\right\}$. The coefficient of $l_{s}, s \neq t$, in (28) is

$$
\begin{equation*}
\sum_{a \in \mathbb{Z}, \eta_{s}^{(\alpha)=1, \eta_{t}(\alpha)=m-1}} 1 \tag{29}
\end{equation*}
$$

which is a count of the total number of nondecreasing sequences $\alpha$ of length $m$ which have the property that rng a contains precisely 1 integer in the interval $\left[P_{s-1}+1, P_{s}\right]$ and $m-1$ integers in $\left[P_{t-1}+1, P_{t}\right]$. Since there are precisely

$$
\begin{equation*}
\binom{p_{t}+m-2}{m-1} \cdot p_{s} \tag{30}
\end{equation*}
$$

such sequences, the value of (29) is (30). Thus (28) for the choice $k_{t}=1, k_{j}=0$ for $j \neq t$ is

$$
l_{t} m\binom{p_{t}+m-1}{m}+\sum_{s \neq t}^{r} l_{s} p_{s}\binom{p_{t}+m-2}{m-1}=0
$$

or

$$
\begin{equation*}
l_{t}\left(p_{t}+m-1\right)+\sum_{s \neq t}^{r} l_{s} p_{s}=0 \tag{31}
\end{equation*}
$$

Consider the system of homogeneous linear equations for $l_{1}, \cdots, l_{r}$ obtained by setting $t=1,2, \cdots, r$ in (31). By setting $k_{1}=\cdots=k_{r}=1$ in (28) we obtain the following additional condition on the $l$ 's:

$$
\begin{equation*}
\sum_{t=1}^{r} l_{t}\left(\sum_{a \in \overline{\mathbf{\Delta}}} \eta_{t}(\alpha)\right)=0 \tag{32}
\end{equation*}
$$

Now by (26)

$$
\sum_{a \in \bar{\Delta}} \eta_{t}(\alpha)=\sum_{\alpha \in \bar{\Delta}} \sum_{j=P_{t-1}+1}^{P_{t}} m_{j}(\alpha)=\sum_{j=P_{t-1}+1}^{P_{t}} q_{j}=p_{t} \cdot q
$$

Thus (32) becomes

$$
\begin{equation*}
\sum_{s=1}^{r} l_{s} p_{s}=0 . \tag{33}
\end{equation*}
$$

Combining the system (31) with (33) we have ( $m-1$ ) $l_{t}=0, t=1, \cdots, r$, and thus $l_{t}=0, t=1, \cdots, r$, contradicting the fact that the $l_{1}, \cdots, l_{r}$ are distinct. Hence $r=1$. In other words, the condition

$$
\begin{equation*}
\operatorname{tr} \Omega_{1, m-1}(A, X)=0, \quad X \in \mathscr{U} \tag{34}
\end{equation*}
$$

implies that all the eigenvalues of $A$ are equal. If we set $X=I_{n}$ in (34) we then see that

$$
\begin{aligned}
\operatorname{tr} \Omega_{1, m-1}\left(A, I_{n}\right) & =\sum_{\alpha \in \bar{\Sigma}} \sum_{t=1}^{m} \lambda_{\alpha(t)}=\sum_{\alpha \in \bar{\Xi}} \sum_{t=1}^{n} m_{t}(\alpha) \lambda_{t} \\
& =\sum_{t=1}^{n} q \cdot \lambda_{t}=q \operatorname{tr} A=0 .
\end{aligned}
$$

Thus we conclude that if $\operatorname{tr} \Omega_{1, m-1}(A, X)=0, X \in \mathfrak{U}$, then all the eigenvalues of $A$ are zero. By repeating precisely the same argument as we gave in Lemma 3, we can conclude that the eigenvalues of both of the normal matrices $A \pm A^{*}$ are zero and hence $A=0$. This completes the proof of Theorem 2.

The proof of Corollary 1 is now obvious, since the conjugate transpose of a generalized doubly stochastic matrix is a matrix of the same kind.

Proof of Theorem 3. We are assuming that $\mathfrak{T}: \operatorname{Hom}(v, v) \rightarrow \operatorname{Hom}(V, v)$ satisfies

$$
\begin{equation*}
\mu_{K}(\mathscr{T}(X))=\mu_{K}(X), \quad X \in \operatorname{Hom}(V, V) \tag{35}
\end{equation*}
$$

Since $\mu_{K}(\xi X)=\xi^{m} \mu_{K}(X)$ it is clear that we may assume $\mathscr{T}\left(I_{n}\right)=I_{n^{*}}$. From (35) we have

$$
\operatorname{tr} K\left(\mathscr{T}\left(I_{n}+x X\right)\right)=\mu_{K}\left(\mathcal{T}\left(I_{n}+x X\right)\right)=\mu_{K}\left(I_{n}+x X\right)=\operatorname{tr} K\left(I_{n}+x X\right)
$$

From (7) with $T=I_{n^{\prime}}, S=X, x_{1}=1, x_{2}=x$ we have

$$
\sum_{r=0}^{m} x^{r} \operatorname{tr} \Omega_{m-r, r}\left(I_{n}, \mathcal{T}(X)\right)=\sum_{r=0}^{m} x^{r} \operatorname{tr} \Omega_{m-r, r}\left(I_{n}, X\right)
$$

and thus

$$
\begin{equation*}
\operatorname{tr} \Omega_{m-r, r}\left(I_{n}, \mathscr{T}(X)\right)=\operatorname{tr} \Omega_{m-r, r}\left(I_{n}, X\right) \tag{36}
\end{equation*}
$$

Let $\kappa_{1}, \cdots, \kappa_{n}$ be the eigenvalues of $X$. Then as we know from (12), the eigenvalues of $\Omega_{m-r, r}\left(I_{n}, X\right)$ are

$$
\begin{equation*}
\sum^{\prime} \prod_{j=1}^{r} \kappa_{a y(j)}, \quad a \in \pi \tag{37}
\end{equation*}
$$

where we recall that the prime indicates that the summation is over all $\gamma \in Q_{r, m}{ }^{*}$ But the expression (37) is precisely the $\pi$ th elementary symmetric function of $\kappa_{a(1)}, \cdots, \kappa_{a(m)}, E_{r}\left(\kappa_{a(1)}, \cdots, \kappa_{a(m)}\right)$. Thus

$$
\begin{equation*}
\operatorname{tr} \Omega_{m-r, r}\left(I_{n^{\prime}} X\right)=\sum_{a \in G_{m, n}} E_{r}\left(\kappa_{\alpha(1)}, \cdots, \kappa_{\alpha(m)}\right) \tag{38}
\end{equation*}
$$

and an elementary induction argument shows that the right-hand side of (38) is precisely

$$
\begin{equation*}
\binom{m+n-1}{m-r} b_{r}\left(\kappa_{1}, \cdots, \kappa_{n}\right) \tag{39}
\end{equation*}
$$

where $h_{r}$ denotes the $r$ th completely symmetric polynomial in the $\kappa$ 's. From (36) and (39) it follows that

$$
b_{r}\left(\kappa_{1}^{\prime}, \cdots, \kappa_{n}^{\prime}\right)=b_{r}\left(\kappa_{1}, \cdots, \kappa_{n}\right), \quad r=1,2, \cdots, m,
$$

where the $\kappa_{1}^{\prime}, \cdots, \kappa_{n}^{\prime}$ are the eigenvalues of $\mathscr{T}(X)$. But by Wronski's relations [13] we know that the completely symmetric polynomials $b_{1}, \cdots, b_{m}$ form an integral polynomial basis for the space of all integral homogeneous symmetric polynomials of degree $m$. Thus it follows that

$$
E_{r}\left(\kappa_{1}^{\prime}, \cdots, \kappa_{n}^{\prime}\right)=E_{r}\left(\kappa_{1}, \cdots, \kappa_{n}\right), \quad r=1,2, \cdots, m
$$

In other words, we have proved that the $r$ th elementary symmetric function of the eigenvalues of both $X$ and $\mathscr{T}(X)$ are equal, $r=1,2, \cdots, m$; and we are in a position to apply a theorem of Marcus and Purves [10] (recently extended by Beasley [2]) which states that any such transformation must have one of the two forms indicated in (3) or (4). This completes the proof of Theorem 3.

The proof of Corollary 2 is an immediate consequence of Theorem 3. For, the eigenvalues of $K(X)$ when $H=S_{m^{\prime}} X \equiv 1, P=V^{(m)}$ are the $\left.\binom{n+m}{m}^{1}\right)$
homogeneous products

$$
\prod_{t=1}^{n} \lambda_{t}^{m_{i}^{(\alpha)}}, \quad \alpha \in G_{m, n}
$$

and hence

$$
\mu_{K}(X)=\operatorname{tr} K(X)=\sum_{a \in G_{m, n}} \prod_{t=1}^{n} \lambda_{t}^{m_{t}^{(a)}}=b_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right) .
$$

Proof of Theorem 4. Once again we can assume that $\mathfrak{T}\left(I_{n}\right)=I_{n}$. Our first task is to determine the structure of the $\bar{\Delta}$ set for $H=A_{m}$ and $\chi \equiv 1$. There are two cases to consider: $m \leq n$ and $m>n$. Let $\omega \in \Gamma_{m, n}$. If the integers in rng $\omega$ are distinct, then let $\alpha$ be the sequence such that $\mathrm{mg} \alpha=\mathrm{mg} \omega$ and $\alpha(1)<\alpha(2)<$ $\ldots<\alpha(m)$, and let $\beta$ be the sequence such that $\mathrm{mg} \beta=\mathrm{mg} \omega$ and $\beta(1)<\beta(2)<$ $\cdots<\beta(m-2)<\beta(m)<\beta(m-1)$. It is clear that $\alpha$ and $\beta$ lie in distinct $A_{m}$-orbits and that any other sequence $\gamma \in \Gamma_{m, n}$ for which mg $\gamma=\mathrm{mg} \omega$ is in the same $A_{m}$-orbit with either $\alpha$ or $\beta$. Thus each such sequence $\omega \in \Gamma_{m, n}$ gives rise to two elements in $\bar{\Delta}$, namely $\alpha$ and $\beta$. On the other hand, if $\omega \in \Gamma_{m_{\rho} n}$ satisfies $m_{t}(\omega) \geq 2$ for some $1 \leq t \leq n$, then it is obvious that there exists a sequence $a$ in $G_{m, n}$ in the same $A_{m}$-orbit with $\omega$. Thus in the case $m \leq n$, the system of distinct representatives may be chosen to be $\bar{\Delta}=G_{m, n} \cup Q_{m, n}^{\prime}$, where $Q_{m, n}^{\prime}$ consists of precisely those sequences $\beta$ for which $\beta(1)<\beta(2)<\cdots<\beta(m-2)<$ $\beta(m)<\beta(m-1)$. In the case that $m>n$, then it is clear that any sequence $\omega \epsilon$ $\bar{\Delta}$ lies in the same $A_{m}$-orbit with a sequence $a \in G_{m, n}$, so that we can choose $\bar{\Delta}=G_{m, n}$. In order to deal with both cases at once we will let $Q_{m, n}^{\prime}=\varnothing$ if $m>n$.

Precisely as in (36) we see that if $\mathfrak{T} \in \mathscr{L}_{1}\left(\mathrm{GL}_{n}(V), K\right)$, then

$$
\begin{equation*}
\operatorname{tr} \Omega_{m-r, r}\left(I_{n}, \mathcal{T}(X)\right)=\operatorname{tr} \Omega_{m-r, r}\left(I_{n}, X\right), \quad X \in M_{n}(\mathbf{C}) \tag{40}
\end{equation*}
$$

The eigenvalues of $\Omega_{m-r, r}\left(I_{n}, X\right)$ are precisely the numbers

$$
E_{r}\left(\kappa_{\alpha(1)}, \cdots, \kappa_{\alpha(m)}\right), \quad \alpha \in \pi,
$$

and thus

$$
\begin{align*}
\operatorname{tr} \Omega_{m-r, r}\left(I_{n}, X\right) & =\sum_{\alpha \in \Lambda} E_{r}\left(\kappa_{\alpha(1)}, \cdots, \kappa_{\alpha(m)}\right) \\
& =\sum_{a \in G_{m, n}} E_{r}\left(\kappa_{\alpha(1)}, \cdots, \kappa_{\alpha(m)}\right)+\sum_{a \in Q_{m, n}^{\prime}} E_{r}\left(\kappa_{\alpha(1)}, \cdots, \kappa_{\alpha(m)}\right)  \tag{41}\\
& =\binom{n+m-1}{m-r} h_{r}\left(\kappa_{1}, \cdots, \kappa_{n}\right)+\binom{n-r}{m-r} E_{r}\left(\kappa_{1}, \cdots, \kappa_{n}\right) .
\end{align*}
$$

In case $m>n$, the elementary symmetric function does not appear in (41) and it is clear from (40) that

$$
\begin{aligned}
b_{r}(\mathcal{T}(X)) & =\binom{n+m-1}{m-r}^{-1} \operatorname{tr} \Omega_{m-r, r}\left(I_{n^{\prime}}, \mathcal{T}(X)\right) \\
& =\binom{n+m-1}{m-r}^{-1} \operatorname{tr} \Omega_{m-r, r}\left(I_{n^{\prime}} X\right) \\
& =b_{r}(X), \quad r=1,2, \cdots, m ; X \in M_{n}(\mathbf{C})
\end{aligned}
$$

Hence by Corollary $2, \mathfrak{T}$ has the required form. On the other hand, if $3 \leq m \leq n$, then using Wronski's relations we have

$$
\begin{aligned}
\operatorname{tr} \Omega_{m-1,1}\left(I_{n^{\prime}} X\right) & =a E_{1}(X), \\
\operatorname{tr} \Omega_{m-2,2}\left(I_{n^{\prime}} X\right) & =b b_{2}(X)+c E_{2}(X) \\
& =b\left(E_{1}^{2}(X)-E_{2}(X)\right)+c E_{2}(X) \\
& =(c-b) E_{2}(X)+b E_{1}^{2}(X), \\
\operatorname{tr} \Omega_{m-3,3}\left(I_{n^{\prime}}, X\right) & =d b_{3}(X)+e E_{3}(X) \\
& =d\left(E_{1}^{3}(X)-2 E_{1}(X) E_{2}(X)+E_{3}(X)\right)+e E_{3}(X) \\
& =(d+e) E_{3}(X)-2 d E_{1}(X) E_{2}(X)+d E_{1}^{3}(X),
\end{aligned}
$$

where $a=\binom{m+n-1}{m}+\binom{n-1}{m-1}, b=\binom{n+m-1}{m}, c=\binom{n-2}{m-2}, d=\binom{n+m-1}{m}$ and $e=\left(\begin{array}{c}n \\ m\end{array}-3 \begin{array}{l}3\end{array}\right)$. Observe that $d+e>0, c-b \neq 0, a \neq 0$; thus the relations (42) allow us to express $E_{3}(X)$ as a polynomial in $\operatorname{tr} \Omega_{m-r, r}\left(I_{n^{\prime}} X\right), r=1,2,3$. It follows from (40), then, that $E_{3}(\mathcal{T}(X))=E_{3}(X), X \in M_{n}^{m-r_{r}}(\mathrm{C})$ and hence we can conclude as before that $\mathfrak{T}$ has the required form. This completes the proof of Theorem 4.

By similar arguments Corollary 3 follows from Theorem 4. For, the eigenvalues of $K(X)$ when $H=A_{m}, \chi \equiv 1$ are the numbers

$$
\prod_{t=1}^{n} \lambda^{m t^{(a)}}, \quad \alpha \in G_{m, n} \cup Q_{m, n}^{\prime}
$$

and hence

$$
\begin{aligned}
\mu_{K}(X) & =\operatorname{tr} K(X)=\sum_{a \in G_{m, n}} \prod_{t=1}^{n} \lambda_{t}^{m_{t}^{(\alpha)}}+\sum_{a \in Q_{m, n}} \prod_{t=1}^{n} \lambda_{t}^{m^{m}}{ }^{(\alpha)} \\
& =k_{m}\left(\lambda_{1}, \cdots, \lambda_{n}\right)
\end{aligned}
$$

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