

Pacific Journal of Mathematics

TRANSFORMATIONS OF SYMMETRIC TENSORS

LARRY JEAN CUMMINGS

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This paper is about linear transformations of the k -fold symmetric tensor product of an n -dimensional vector space V which carry nonzero decomposable tensors to nonzero decomposable tensors. The main theorem shows that every such transformation is induced by a nonsingular transformation of V provided both

- (i) the field has characteristic either 0 or a prime greater than k and every polynomial over the field with degree at n is a product of linear factors.
- (ii) $n > k + 1$.

Condition (i) includes the important special case where the field is algebraically closed with characteristic 0.

The linear transformations which preserve decomposable tensors in the skew-symmetric case have been studied in two papers by Westwick [6, 8]. In [6] he showed that if the field is algebraically closed then the transformation is induced by a linear transformation of V except, possibly, when the dimension of V is $2k$. In the latter case the transformation may be the composition of one induced by a linear transformation of V and one induced by a correlation of the k -dimensional subspaces of V . A series of papers [3, 4, 7, 2] has been devoted to linear transformations which preserve decomposable tensors in the case of the full tensor product.

Our result partially answers a question first raised by Marcus and Newman in [5]. They asked for necessary and sufficient conditions in order that every decomposable mapping of the space of k -fold symmetric tensors be induced.

1. Preliminaries. Let V^k denote the k -fold Cartesian product of V where $k > 1$. A k -fold symmetric tensor space (or rank k symmetric tensor space) is a vector space denoted by $\mathbf{V}_k V$ together with a fixed multilinear symmetric mapping $\sigma: V^k \rightarrow \mathbf{V}_k V$ which is universal for multilinear and symmetric mappings of $\mathbf{V}_k V$. We assume that $\mathbf{V}_k V$ is generated by the image of σ . Thus, if W is any vector space and $g: V^k \rightarrow W$ is both multilinear and symmetric then g has a unique extension $h: \mathbf{V}_k V \rightarrow W$ such that

$$(1.1) \quad \begin{array}{ccc} & W & \\ g \nearrow & & \nwarrow h \\ V^k & \xrightarrow{\sigma} & \mathbf{V}_k V \end{array}$$

is commutative and $\mathbf{V}_k V$ is isomorphic to any other vector space with this property. In particular, if $A: V \rightarrow V$ is linear then the assignment

$$(x_1, \dots, x_k) \longmapsto Ax_1 \vee \dots \vee Ax_k$$

is a multilinear and symmetric mapping of V^k . We will denote its unique linear extension to $\mathbf{V}_k V$ by $\mathbf{V}_k A$.

The *decomposable symmetric tensors* or "symmetric products" are images under σ of k -tuples in V^k . For convenience we denote $\sigma(x_1, \dots, x_k)$ by $x_1 \vee \dots \vee x_k$. A subspace s of $\mathbf{V}_k V$ is decomposable if $S \subseteq \sigma(V^k)$. *Trivial decomposable subspaces* are the zero subspace and the 1-dimensional subspaces whose elements are scalar multiples of a single nonzero decomposable symmetric tensor. If V and F satisfy (i) and (ii) the maximal decomposable subspaces of $\mathbf{V}_k V$ were determined in [1].

A symmetric product is zero if and only if at least one of its factors is zero. More generally, if

$$x_1 \vee \dots \vee x_k = y_1 \vee \dots \vee y_k \neq 0$$

then there are scalars $\lambda_1, \dots, \lambda_k$ such that $\lambda_1 \dots \lambda_k = 1$ and

$$(1.2) \quad x_i = \lambda_i y_{\pi(i)} \quad i = 1, \dots, k.$$

Here $\pi \in S_k$, the symmetric group on $\{1, \dots, k\}$.

A linear transformation $f: \mathbf{V}_k V \rightarrow \mathbf{V}_k V$ is *decomposable* if

$$f(\sigma(V^k)) \subseteq \sigma(V^k)$$

and

$$(1.3) \quad \ker f \cap \sigma(V^k) = 0.$$

If V is an n -dimensional vector space then the dimension of $\mathbf{V}_k V$ is $\binom{n+k-1}{k}$.

2. Type 1 subspaces and associate mappings. Subspaces in $\mathbf{V}_k V$ of the form

$$(2.1) \quad M = x_1 \vee \dots \vee x_{k-1} \vee V \quad k > 1$$

where x_1, \dots, x_{k-1} are fixed nonzero vectors in V are always decomposable because of the multilinearity of the mapping σ . It is convenient to call these *type 1 subspaces*. The 1-dimensional subspaces $\langle x_1 \rangle, \dots, \langle x_{k-1} \rangle$ are called the *factors* of M .

PROPOSITION 1. *If F is a field whose characteristic (if any) is*

not less than k then

$$(2.2) \quad x_1 \vee \cdots \vee x_{k-1} \vee V = x'_1 \vee \cdots \vee x'_{k-1} \vee V \quad \text{in } \mathbf{V}_k V$$

implies

$$\langle x_1 \vee \cdots \vee x_{k-1} \rangle = \langle x'_1 \vee \cdots \vee x'_{k-1} \rangle \quad \text{in } \mathbf{V}_{k-1} V.$$

Proof. This proof requires the choice of a vector not in the set-theoretic union

$$(2.3) \quad \langle x_1 \rangle \cup \cdots \cup \langle x_{k-1} \rangle.$$

By Lemma 12 of [1, p. 73] we know that if V were the union (2.3) then the cardinality of F could not exceed the finite integer $k-1$. This would mean that the characteristic of F exceeds the cardinality of F . Accordingly we may choose v in V not in the union (2.3) and (2.2) implies the existence of a u in V satisfying

$$x_1 \vee \cdots \vee x_{k-1} \vee u = x'_1 \vee \cdots \vee x'_{k-1} \vee v.$$

By the choice of v and (1.2) there is a nonzero scalar λ for which $u = \lambda v$ and

$$x_i = \lambda_i x'_{\pi(i)} \quad i = 1, \dots, k-1$$

where $\pi \in S_{k-1}$ and $1 = \lambda \prod \lambda_i$. Therefore,

$$\lambda x_1 \vee \cdots \vee x_{k-1} = x'_1 \vee \cdots \vee x'_{k-1}$$

in $\mathbf{V}_{k-1} V$.

Hereafter we will assume that F satisfies the hypothesis of Proposition 1.

A *type 1 mapping* is a decomposable mapping of $\mathbf{V}_k V$ for which the image of every type 1 subspace is again a type 1 subspace. If f is a type 1 mapping and M is the type 1 subspace (2.1) then we may choose nonzero vectors y_1, \dots, y_{k-1} in V such that

$$(2.4) \quad f(M) = y_1 \vee \cdots \vee y_{k-1} \vee V.$$

We obtain a well-defined linear mapping A of V by setting $Au = v$ if

$$(2.5) \quad f(x_1 \vee \cdots \vee x_{k-1} \vee u) = y_1 \vee \cdots \vee y_{k-1} \vee v.$$

The mapping A will be called an *associate mapping* of f with respect to M . In general, the associate map defined by (2.5) depends not only on M and f but the choice of the vectors y_1, \dots, y_{k-1} as well.

PROPOSITION 2. *Any two associate mappings of a type 1 mapping with respect to the same type 1 subspace are multiples.*

Proof. This follows easily from Proposition 1 and (1.1).

PROPOSITION 3. *Every associate of a type 1 mapping is non-singular.*

Proof. Let A be an associate of a type 1 mapping f with respect to (2.1) and suppose $A(u) = A(u')$ for some vectors u, u' in V . From (2.5) we have

$$f(x_1 \vee \cdots \vee x_{k-1} \vee u) = f(x_1 \vee \cdots \vee x_{k-1} \vee u').$$

Since f is linear and decomposable we have

$$x_1 \vee \cdots \vee x_{k-1} \vee (u - u') = 0$$

which implies $u = u'$.

Two type 1 subspaces will be called *adjacent* if they have exactly $k - 2$ common factors (counting multiplicity). Accordingly a typical pair of adjacent subspaces may be written in the form

$$(2.6) \quad M_i = x_1 \vee \cdots \vee x_{k-1} \vee z_i \vee V \quad i = 1, 2$$

where z_1, z_2 are two independent vectors of V and x_1, \dots, x_{k-1} are arbitrary nonzero vectors.

Two arbitrary type 1 subspaces are always connected by a chain of adjacent subspaces; explicitly, if

$$(2.7) \quad M = x_1 \vee \cdots \vee x_{k-1} \vee V$$

and

$$N = y_1 \vee \cdots \vee y_{k-1} \vee V$$

then M_p is adjacent to M_{p+1} where

$$(2.8) \quad M_p = x_1 \vee \cdots \vee x_{k-p-1} \vee y_1 \vee \cdots \vee y_p \vee V \quad p = 1, \dots, k - 2$$

and we take $M = M_0$ and $N = M_{k-1}$.

PROPOSITION 4. *Two type 1 subspaces M and N are adjacent if and only if $\dim M \cap N = 1$. Otherwise $M \cap N = 0$ whenever M and N are distinct.*

Proof. Consider the adjacent type 1 subspaces (2.6). If $t \in M_1 \cap M_2$ then there exist vectors u and v in V such that

$$(2.9) \quad t = x_1 \vee \cdots \vee x_{k-2} \vee z_1 \vee u = x_1 \vee \cdots \vee x_{k-2} \vee z_2 \vee v.$$

Now the multilinear and symmetric mapping $g_p(x): V^p \rightarrow \mathbf{V}_{p+1}V$ defined for each $p = 2, \dots, k-1$ by

$$(2.10) \quad (v_1, \dots, v_p) \longmapsto x \vee v_1 \vee \cdots \vee v_p$$

extends as in (1.1) to a linear mapping $h_p(x): \mathbf{V}_p V \rightarrow \mathbf{V}_{p+1}V$. If the vector x in (2.10) is nonzero then each $h_p(x)$ is injective and so is the composite

$$h = h_{k-1}(x_1) \cdots h_{k-i}(x_i) \cdots h_2(x_{k-2}).$$

Thus (2.9) is just

$$h(z_1 \vee u) = h(z_2 \vee v)$$

and so

$$z_1 \vee u = z_2 \vee v.$$

Since z_1 and z_2 are independent (1.2) implies that u is a scalar multiple of z_2 . Therefore

$$(2.11) \quad M_1 \cap M_2 = \langle x_1 \vee \cdots \vee x_{k-2} \vee z_1 \vee z_2 \rangle.$$

Now consider an arbitrary pair of type 1 subspaces (2.7) and suppose they have nonzero intersection. Let

$$t = x_1 \vee \cdots \vee x_{k-1} \vee u = y_1 \vee \cdots \vee y_{k-1} \vee v$$

be a nonzero element of the intersection. If $\langle u \rangle = \langle v \rangle$ then by (1.2) we have $M_1 = M_2$ and otherwise M_1 and M_2 must have exactly $k-2$ common factors.

PROPOSITION 5. *The images of adjacent type 1 subspaces under type 1 mappings are adjacent provided the underlying field satisfies (i).*

Proof. Consider the adjacent type 1 subspaces (2.6). We know from Proposition 4 that

$$M_1 \cap M_2 = \langle x_1 \vee \cdots \vee x_{k-2} \vee z_1 \vee z_2 \rangle.$$

If f is a type 1 mapping then $f(M_1) \cap f(M_2)$ is nonzero and Proposition 4 yields the desired conclusion provided $f(M_1)$ and $f(M_2)$ are distinct. We complete the proof by showing that the images of adjacent subspaces are always distinct.

Consider the two linear mappings $A_i: V \rightarrow \mathbf{V}_k V$ defined by

$$A_i(v) = f(x_1 \vee \cdots \vee x_{k-1} \vee z_i \vee v) \quad i = 1, 2.$$

It follows that they are injective because f is linear and decomposable. Suppose $\text{range } A_1 = \text{range } A_2$ and let $A_2^{-1}: \text{range } A_2 \rightarrow V$ be the inverse of A_2 . Then $A_2^{-1}A_1$ is a well-defined linear transformation of V . Because of (i), $A_2^{-1}A_1$ has at least one characteristic value, say λ . If u is a corresponding characteristic vector then $A_1u = \lambda A_2u$. That is,

$$f(x_1 \vee \cdots \vee x_{k-1} \vee z_1 \vee u) = \lambda f(x_1 \vee \cdots \vee x_{k-1} \vee z_2 \vee u).$$

Since f is linear and decomposable we obtain $z_1 = \lambda z_2$, contradicting the assumption that M_1 and M_2 are adjacent.

Any collection of two or more type 1 subspaces in $\bigvee_k V$ ($k > 2$) will be called an *adjacent family* if there are vectors x_1, \dots, x_{k-2} in V such that any subspace in the collection can be written as

$$x_1 \vee \cdots \vee x_{k-2} \vee u \vee V$$

for some vector $u \in V$. When $k = 2$ any collection containing at least two distinct type 1 subspaces will be called an adjacent family. Of course every pair of adjacent type 1 subspaces constitutes an adjacent family, but a collection of three or more need not be, as is easily seen by example.

PROPOSITION 6. *Any collection of more than k pair-wise adjacent type 1 subspaces in $\bigvee_k V$ is an adjacent family.*

Proof. We assign to each type 1 subspace (2.1) the set

$$\{(\langle x_i \rangle, i) \mid i = 1, \dots, k-1\}$$

which always contains $k-1$ distinct elements even if (2.1) does not have distinct factors.

The proposition now follows from the combinatorial result that a collection of more than k finite sets each containing $k-1$ elements which intersect pair-wise in $k-2$ elements always intersect in the same set of $k-2$ elements:

If $k = 2$ there is nothing to prove. If $k > 2$ let X and Y be any two sets of the collection. There are elements a and b such that

$$X = (X \cap Y) \cup \{a\}$$

and

$$Y = (X \cap Y) \cup \{b\}.$$

Because any two sets in the collection intersect in $k - 2$ elements, any set of the collection not containing $X \cap Y$ must contain both a and b and intersect $X \cap Y$ in exactly $k - 3$ elements. But there are at most $k - 2 = \binom{k-2}{k-3}$ distinct such sets. Therefore, the collection must contain at least one set Z distinct from X and Y but which contains $X \cap Y$. Let

$$Z = X \cap Y \cup \{c\}$$

and suppose there exists a set W in the collection not containing $X \cap Y$. Then $\{a, b, c\} \subseteq W$, contradicting the hypothesis that $X \cap W$ has $k - 2$ elements.

3. Main results. A collection of vectors in an n -dimensional vector space is said to be in *general position* when any n vectors chosen from the collection form a basis of V . The following well known lemma about vectors in general position will be used in showing that any two associate mappings of a type 1 mapping are multiples whenever $n > 2$ and the underlying field is infinite.

LEMMA 1. *If $m \geq n$ then an n -dimensional vector space over an infinite field always contains m vectors in general position.*

LEMMA 2. *Let z_1, \dots, z_m be any finite set of vectors in an n -dimensional vector space over an infinite field. If $A: V \rightarrow V$ is nonsingular and B is any other linear mapping of V satisfying*

$$(3.1) \quad \langle A(x) \rangle = \langle B(x) \rangle$$

for all vectors x not in $S = \langle z_1 \rangle \cup \dots \cup \langle z_m \rangle$ then there is a scalar λ such that $B = \lambda A$.

Proof. Since F is infinite Lemma 12 of [1] and induction show the existence of a basis of V disjoint from the set S . If b_1, \dots, b_n is such a basis let $\lambda_1, \dots, \lambda_n$ be scalars such that

$$(3.2) \quad B(b_i) = \lambda_i A(b_i) \quad i = 1, \dots, n.$$

Since F is infinite we may choose a vector $v = \sum \alpha_i b_i$ not in S but all of whose coordinates with respect to b_1, \dots, b_n are non-zero. Then (3.1) and (3.2) imply the existence of a scalar λ such that

$$\sum \alpha_i \lambda_i A(b_i) = \sum \lambda \alpha_i A(b_i).$$

Since A is nonsingular we have $\lambda_1 = \lambda_2 = \dots = \lambda_n = \lambda$.

REMARK. In (i) we assume that every polynomial of degree at

most n splits completely over the underlying field. This means that the field is necessarily infinite since the polynomial ring over a finite field has irreducible elements of every degree. Thus Lemmas 1 and 2 are immediately applicable in the following theorems.

THEOREM 1. *The associate mappings of a type 1 mapping of $\mathbb{V}kV$ are a 1-dimensional subspace of the linear mappings of V , provided $\dim V > 2$ and F satisfies (i).*

Proof. We show first that an associate map of a type 1 mapping f with respect to one of type 1 subspaces (2.6) is always a scalar multiple of every associate mapping of the other. By Lemma 1 we complete the vectors z_1, z_2 to a set z_1, \dots, z_m in general position where $m = \text{Max}\{k, \dim V\}$. As in the proof of the Proposition 1 we may choose a vector z_{m+1} not in the set-theoretic union $\langle z_1 \rangle \cup \dots \cup \langle z_m \rangle$. Then the subspaces

$$M_i = x_1 \vee \dots \vee x_{k-2} \vee z_i \vee V \quad i = 1, \dots, m+1$$

are an adjacent family. The images of these subspaces form a family of pair-wise adjacent subspaces by Proposition 5. They form an adjacent family by Proposition 6 and the choice of m . Thus we may choose vectors $y_1, \dots, y_{k-2}; w_1, \dots, w_{m+1}$ in V such that

$$(3.3) \quad f(M_i) = y_1 \vee \dots \vee y_{k-2} \vee w_i \vee V \quad i = 1, \dots, m+1.$$

We proceed to examine the effect of f on the intersections $M_i \cap M_{m+1}; i = 1, 2$. By (3.3)

$$\begin{aligned} f(x_1 \vee \dots \vee x_{k-2} \vee z_i \vee z_{m+1}) &= y_1 \vee \dots \vee y_{k-2} \vee w_i \vee A_i(z_{m+1}) \\ &= y_1 \vee \dots \vee y_{k-2} \vee w_{m+1} \vee A_{m+1}(z_i) \\ &\quad i = 1, 2. \end{aligned}$$

where A_i denotes any associate map of M_i under f and A_{m+1} is an associate of M_{m+1} . It follows that $\langle w_{m+1} \rangle = \langle A_i(z_{m+1}) \rangle$ for $i = 1, 2$ because w_{m+1} is not in $\langle w_1 \rangle \cup \langle w_2 \rangle$. Since z_{m+1} is restricted only by its exclusion from $\langle z_1 \rangle \cup \dots \cup \langle z_m \rangle$ Lemma 2 applies and yields a scalar γ such that $A_1 = \gamma A_2$.

To complete the proof we need only consider an arbitrary pair of type 1 subspaces (2.7) and a chain (2.8) of adjacent subspaces between them. If A_p is an associate map of M_p then we have just shown the existence of a scalar γ_p such that

$$A_p = \gamma_p A_{p+1} \quad p = 0, \dots, k-2.$$

Therefore, $A_0 = \gamma_0 \dots \gamma_{k-2} A_{k-1}$.

REMARK. If $\dim V = 1$ then $\mathbf{V}_k V = 1$ and $L(\mathbf{V}_k V, \mathbf{V}_k V) \cong F$. Hence $L(\mathbf{V}_k V, \mathbf{V}_k V)$ consists of induced mappings if and only if every polynomial of the form $x^k - a$ has a root in F .

THEOREM 2. *Every type 1 mapping of $\mathbf{V}_k V$ is induced by an associate mapping, provided $\dim V > 2$ and F satisfies (i).*

Proof. Let $x = x_1 \vee \cdots \vee x_k$ be any nonzero product of $\mathbf{V}_k V$. The trivial subspace $\langle x \rangle$ is the intersection of the k type 1 subspaces

$$(3.4) \quad T_i = x_1 \vee \cdots \vee \hat{x}_i \vee \cdots \vee x_k \vee V \quad i = 1, \dots, k.$$

By Theorem 1 the associate mappings of a type 1 mapping f with respect to the subspaces (3.4) are scalar multiples of one another. If A is any one of them then Theorem 1 and definition (2.5) show then that Ax_i must be a factor of $f(x)$ for each $i = 1, \dots, k$. Thus, if x has distinct factors it follows from (1.2) and Proposition 3 that

$$(3.5) \quad f(x) = \lambda_x Ax_1 \vee \cdots \vee Ax_k$$

for some scalar λ_x and

$$(3.6) \quad f(T_i) = Ax_1 \vee \cdots \vee \widehat{Ax_i} \vee \cdots \vee Ax_k \vee V \quad i = 1, \dots, k.$$

We next verify (3.6) when the factors $\langle x_1 \rangle, \dots, \langle x_k \rangle$ are not necessarily distinct. To this end consider a chain of adjacent subspaces (2.8) where we suppose M_{k-1} has arbitrary factors and take the factors of M_0 as distinct and distinct from the factors of M_{k-1} . This we may always do since any field satisfying (i) must be infinite. (See the remark following Lemma 2.) Thus (3.6) may be applied to M_0 which contains $z_1 = x_1 \vee \cdots \vee x_{k-1} \vee y_1$. By Theorem 1 there is a scalar λ for which

$$(3.7) \quad f(z_1) = \lambda Ax_1 \vee \cdots \vee Ax_{k-1} \vee Ay_1.$$

Therefore the $k - 1$ factors of $f(M_1)$ must be among the factors of (3.7). Now $\langle Ay_1 \rangle$ could not be excluded because then M_0 and M_1 would have the same type 1 subspace as image, contradicting Proposition 5. If, say, Ax_1 were excluded then

$$f(M_1) = Ax_2 \vee \cdots \vee Ax_{k-1} \vee Ay_1 \vee V$$

and Theorem 1 yields

$$(3.8) \quad f(z_1) = \lambda_1 Ax_2 \vee \cdots \vee Ay_1 \vee Ax_{k-1}$$

for some scalar λ_1 .

Comparison of (3.7) and (3.8) shows that Ax_{k-1} would be a scalar

multiple of either Ay_i or some Ax_i with $1 \leq i < k - 1$. Hence

$$f(M_1) = Ax_1 \vee \cdots \vee Ax_{k-2} \vee Ay_1 \vee V.$$

Suppose it has been shown that

$$(3.9) \quad f(M_p) = Ax_1 \vee \cdots \vee Ax_{k-p-1} \vee Ay_1 \vee \cdots \vee Ay_p \vee V$$

for some p , $1 < p \leq k - 2$. Since

$$M_p \cap M_{p+1} = \langle x_1 \vee \cdots \vee x_{k-p-1} \vee y_1 \vee \cdots \vee y_{p+1} \rangle$$

(3.9) implies that $f(M_{p+1})$ contains

$$(3.10) \quad Ax_1 \vee \cdots \vee Ax_{k-p-1} \vee Ay_1 \vee \cdots \vee Ay_{p+1}$$

and so the $k - 1$ factors of $f(M_{p+1})$ are among the factors of (3.10). Arguing as before we see that Ay_{p+1} must be a factor of $f(M_{p+1})$ since otherwise the images of $f(M_p)$ and $f(M_{p+1})$ would coincide. If, say, Ax_1 were not a factor then

$$f(M_{p+1}) = Ax_2 \vee \cdots \vee Ax_{k-p-1} \vee Ay_1 \vee \cdots \vee Ay_{p+1} \vee V$$

and by Theorem 1 there is a scalar μ for which

$$(3.11) \quad \begin{aligned} & f(x_1 \vee \cdots \vee x_{k-p-1} \vee y_1 \vee \cdots \vee y_{p+1}) \\ &= \mu Ax_2 \vee \cdots \vee Ax_{k-p-1} \vee Ay_1 \vee \cdots \vee Ay_{p+1} \vee Ax_{k-p-1}. \end{aligned}$$

Comparison of (3.10) and (3.11) shows that Ax_{k-p-1} would be either a multiple of some Ay_i , $1 \leq i \leq p + 1$, or some Ax_j , $1 \leq j < k - p - 1$, contradicting the assumption that the factors of M_0 are distinct and distinct from the factors of M_{k-1} .

Since any product x is in some type 1 subspace we have shown that $f(x) = \lambda_x(\mathbf{V}_k A)(x)$ for some scalar λ_x . If x and y are products in the same type 1 subspace a simple comparison argument shows that $\lambda_x = \lambda_y$. Denote the common value by λ . When x and y are arbitrary products we obtain the same result by considering type 1 subspaces containing them and a chain (2.8) between the subspaces since any two of the latter have 1-dimensional intersections. Because the field always contains a root of $x^k - \lambda = 0$ by (i), we have shown that f is induced by $\lambda^{1/k} A$.

THEOREM 3. *Every decomposable mapping of $\mathbf{V}_k V$ is induced by a nonsingular mapping of V , provided V is a finite dimensional vector space satisfying (i) and (ii).*

Proof. Because of the previous theorem we need only show with the additional hypothesis that every decomposable mapping of

$\mathbf{V}_k V$ is type 1. If M is any type 1 subspace and f decomposable then $f(M)$ is a decomposable subspace and hence contained in a maximal decomposable subspace of $\mathbf{V}_k V$. In [1] the maximal decomposable subspaces of $\mathbf{V}_k V$ were determined for the case when V satisfies the hypothesis of this theorem. The subspaces are

- (a) type 1 subspaces
- (b) type r subspaces which are of the form

$$x_1 \vee \cdots \vee x_{k-r} \vee S \vee \cdots \vee S$$

where $1 < r \leq k$ and S is a 2-dimensional subspace of V .

Those subspaces of type $r > 1$ have dimension $r + 1$. If the maximal decomposable subspace containing $f(M)$ was one of these types then $\dim V \leq r + 1 \leq k + 1$ by (1.3) because every type 1 subspace has the same dimension as V . The hypothesis $\dim V > k + 1$ thus implies that the maximal decomposable subspace containing $f(M)$ is type 1 and therefore f is type 1.

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Received June 10, 1971 and in revised form October 22, 1971. Part of this work was contained in the author's thesis written under R. Westwick at the University of British Columbia. The author is also indebted to conversations with B. N. Moyls.

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PUBLISHED BY PACIFIC JOURNAL OF MATHEMATICS, A NON-PROFIT CORPORATION

Printed at Kokusai Bunken Insatsusha (International Academic Printing Co., Ltd.), 270, 3-chome Totsuka-cho, Shinjuku-ku, Tokyo 160, Japan.

Pacific Journal of Mathematics

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