## A THEOREM ON QUADRATIC FORMS AND ITS APPLICATION IN THE CALCULUS OF VARIATIONS\*

BY

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1. Principal result. Let P(z),  $Q_1(z)$ ,  $\cdots$ ,  $Q_r(z)$  be quadratic forms in the real variables  $z_1, \cdots, z_m$  with real coefficients having the following properties: the form P(z) is positive at each point  $(z) \neq (0)$  at which the forms  $Q_1(z), \cdots, Q_r(z)$  vanish simultaneously; for every set of constants  $S_1, \cdots, S_r$ , not all zero, the quadratic form  $S_iQ_i(z)^{\dagger}$  is indefinite; for every linear subspace L on which the quadratic forms  $Q_1, \cdots, Q_r$  do not vanish simultaneously, save at (z) = (0), there is a linear combination  $S_iQ_i(z)$  of these forms which is positive definite on L. The principal result given in the present note is given in the following:

THEOREM 1. Under the above hypotheses there exists a set of constants  $S_1, \dots, S_r$  such that the quadratic form

$$P(z, S) = P(z) + S_i Q_i(z) \qquad (i = 1, \dots, r)$$

is positive definite in the variables  $z_1, \dots, z_m$ . The last two hypotheses may be dropped when r = 1.

This theorem was proposed in a somewhat different form by Bliss in 1937 in a seminar on the calculus of variations. It is useful in sufficiency proofs for multiple integral problems. Proofs of the theorem for the case r=1 were given in the seminar by Albert,  $\ddagger$  Reid,  $\S$  McShane and Hestenes, each using a different method. The last two of these proofs have not been published. The proof of Theorem 1 here given is due to McShane and is an immediate extension of the one given by McShane for the case r=1. In the next section Theorem 1 will be applied to the case in which  $(z)=(x_1,\cdots,x_n,y_1,\cdots,y_n)$  and the forms  $Q_1,\cdots,Q_r$  are the two-rowed minors of the  $n\times 2$ -dimensional matrix  $(x_k,y_k)$ . If n=2 then r=1 and Theorem 2 below is an immediate consequence of Theorem 1. If n>2 the further result described in Theorem 3 below is

<sup>\*</sup> Presented separately to the Society, by McShane on April 8, 1939, and by M. R. Hestenes on December 29, 1939; received by the editors November 29, 1939.

<sup>†</sup> A repeated index denotes summation.

<sup>‡</sup> A. A. Albert, A quadratic form problem in the calculus of variations, Bulletin of the American Mathematical Society, vol. 44 (1938), pp. 250-252.

<sup>§</sup> W. T. Reid, A theorem on quadratic forms, Bulletin of the American Mathematical Society, vol. 44 (1938), pp. 437-440.

needed. This further result for the case n=3 has been established by Mc-Shane.\* The proofs of Theorems 3, 4 and 5 here given are due to Hestenes.†

The proof of Theorem 1 is based on two lemmas, the first of which is the following:

LEMMA A. Let m(S) be the minimum of the quadratic form P(z, S) on the unit (m-1)-sphere  $z_k z_k = 1$ . There exists a set of constants  $(S_0)$  which maximizes the function m(S).

For a set of constants  $S_1, \dots, S_r$  let M(S) be the maximum of the quadratic form  $-S_iQ_i$  on the unit sphere  $z_kz_k=1$ . We have M(S)>0 if  $(S)\neq (0)$  since the quadratic form  $-S_iQ_i$  is indefinite by hypothesis. Moreover M(hS)=hM(S) for every positive number h. Let  $M_0$  be the minimum of M(S) on the set  $S_iS_i=1$ . Since this minimum is attained, we have  $M_0>0$ . Finally let N be the maximum of P(z) on the set  $z_kz_k=1$ . Then for every set of constants (S) there is a point (z) such that  $z_kz_k=1$ ,  $-S_iQ_i(z)=M(S)\geq hM_0$ ,  $S_iS_i=h^2$  and

$$m(S) \leq P(z, S) \leq N - M(S) \leq N - hM_0$$
.

It follows that there is a positive constant  $h_0$  such that when  $h > h_0$  we have m(S) < m(0) for every set (S) with  $S_iS_i = h^2$ . There is accordingly a set  $(S_0)$  such that the relation  $m(S) \le m(S_0)$  holds for every set (S), as was to be proved.

Lemma B. Let  $(S_0)$  be a set of constants which maximizes the function m(S) and set  $m_0 = m(S_0)$ . The set of points (z) satisfying the equation  $P(z, S_0) = m_0 z_k z_k$  forms a linear space L. There is no quadratic form  $S_i Q_i$  which is positive definite on L.

The set L consists of all points z at which the function  $P(z, S_0) - m_0 z_k z_k$  attains its minimum value 0. Hence the partial derivatives of that function vanish on L, and therefore the equations

$$(1) P_{z_k}(z, S_0) = 2m_0 z_k$$

hold for all z in L. Conversely, if we multiply equations (1) by  $z_k$  and sum on k we see that every point z which satisfies (1) lies in L. Thus L is the set of solutions of the linear equations (1), and is therefore linear. Suppose now that there exists a quadratic form  $Q = S_iQ_i$  that is positive definite on L. Let K be the unit sphere  $z_kz_k = 1$  and  $L_1$  the set of points in L on K. Choose b > 0 such

<sup>\*</sup> E. J. McShane, The condition of Legendre for double integral problems of the calculus of variations, Bulletin of the American Mathematical Society, abstract 45-5-209.

<sup>†</sup> M. R. Hestenes, A theorem on quadratic forms and its application in the calculus of variations, Bulletin of the American Mathematical Society, abstract 46-1-83.

that Q(z) > b on  $L_1$ , and let N be a neighborhood of  $L_1$  relative to K on which Q(z) > b. Let  $m_1$  be the minimum of  $P(z, S_0)$  on the closed set K - N. Then  $m_1 > m_0$ . It follows that for a sufficiently small positive constant h one will have

$$P(z, S_0 + hS) = P(z, S_0) + hQ > m_0$$

on K-N. But  $P(z, S_0+hS) > m_0+hb$  on N, and hence  $m(S_0+hS) > m(S_0)$ , contrary to our choice of the set  $(S_0)$ . This proves Lemma B.

In order to prove Theorem 1 we note that there is a point  $(z) \neq (0)$  on the set L described above at which the forms  $Q_1, \dots, Q_r$  vanish simultaneously. Otherwise by virtue of the last hypothesis made in Theorem 1 there would exist a quadratic form  $S_iQ_i(z)$  which is positive definite on L, contrary to Lemma B. At this point (z) we have accordingly  $m_0z_kz_k=P(z,S_0)=P(z)>0$  in view of the first hypothesis of the theorem. It follows that  $m_0>0$  and hence that  $P(z,S_0)$  is positive definite. This proves the first statement in the theorem. The second statement is readily verified. It should be observed that the hypotheses of Theorem 1 imply that there is a point  $(z) \neq (0)$  at which the forms  $Q_1, \dots, Q_r$  vanish simultaneously. Otherwise the last hypothesis would imply the existence of a positive definite form  $S_iQ_i$ , contrary to the second hypothesis.

The last sentence in Theorem 1 remains to be established. It is easy to see that if r=1 the last hypothesis in §1 is automatically satisfied. Suppose then that r=1 and that Q is not indefinite. Let K be the sphere  $z_k z_k = 1$ . If  $Q_1$  is positive definite, the sum P(z, S) is positive on K, hence is a positive definite form, provided that  $S_1$  is large. If  $Q_1$  is negative definite, P(z, S) is positive definite provided that  $-S_1$  is large. If  $Q_1$  is semi-definite, say positive, let L be the set on which  $Q_1$  vanishes. As in Lemma B, this is linear. On LK the form P(z) is positive; it then remains positive on a neighborhood N of LK relative to K. On K-N the form  $Q_1$  is positive. Choose  $S_1$  large enough so that P(z, S) is positive on K-N. On N we have  $P(z, S) \ge P(z) > 0$ , so that P(z, S) is positive on K, and is therefore a positive definite form.

In the proof of the first part of Theorem 1 made above we have established essentially the following more general result:

COROLLARY. Suppose the last two of the hypotheses made in Theorem 1 hold, and let m be the minimum of the form P(z) on the set of points (z) satisfying the conditions  $Q_1(z) = \cdots = Q_r(z) = 0$ ,  $z_k z_k = 1$ . There exists a set of constants  $S_i$  such that the inequality  $P(z, S) \ge m z_k z_k$  holds for all points (z).

2. A further result. Consider now the case in which the space of points (z) described in the last section is of dimension 2n. For the purposes of this sec-

tion it will be convenient to denote the ith  $(i \le n)$  and the (n+i)th coordinates of (z) by  $x_i$  and  $y_i$  respectively. Thus the points of our space will be denoted by the symbol (x, y). To each point (x, y) there is associated a  $n \times 2$ -dimensional matrix  $(x_i, y_i)$  whose ith row is composed of the coordinates  $x_i, y_i$  of the point (x, y). This matrix will be used below to classify the points (x, y) of our space. By a quadratic form in the variables (x, y) will be meant an expression of the form

$$R_{ik}x_ix_k + S_{ik}x_iy_k + T_{ik}y_iy_k$$
.

In particular the expression  $S_{ik}x_iy_k$  is a quadratic form in the variables (x, y). Finally by a linear space L of points (x, y) is meant a subspace such that if (x, y) and (x', y') belong to L so also does the point (ax+bx', ay+by'), where a and b are arbitrary real constants.

The results described in the last section will be used to prove the following:

THEOREM 2. Let P(x, y) be a quadratic form in the 2n real variables  $x_1, \dots, x_n, y_1, \dots, y_n$  with real coefficients. Suppose that the inequality P(x, y) > 0 holds whenever the  $n \times 2$ -dimensional matrix  $(x_i, y_i)$  has rank 1. Then there exists an n-rowed skew-symmetric matrix  $S = (S_{ik})$  such that the quadratic form  $P(x, y) + S_{ik}x_iy_k$   $(i, k = 1, \dots, n)$  is positive definite.

To prove this result let  $Q_{ik}$  (i < k) be the quadratic form  $x_i y_k - x_k y_i$ . A linear combination  $S_{ik}Q_{ik}$  (summed with i < k) is easily seen to be equal to  $S_{ik}x_i y_k$  (summed for all i, k) if we set  $S_{ii} = 0, S_{ki} = -S_{ik}$  (i < k). Thus we see that the theorem will be established if we show that the hypotheses of Theorem 1 with  $Q_i$  replaced by  $Q_{ik}$  are satisfied. The first hypothesis holds since the matrix  $(x_i, y_i) \neq (0, 0)$  has rank 1 if and only if the forms  $Q_{ik}$  vanish simultaneously. Moreover a linear combination  $S_{ik}Q_{ik}$  with  $S \neq 0$  is indefinite. Finally the last hypothesis of Theorem 1 holds by virtue of the following:

THEOREM 3. Let L be a linear set of points (x, y) such that the  $n \times 2$ -dimensional matrix  $(x_i, y_i)$  has rank 2 at each point  $(x, y) \neq (0, 0)$  on L. There exists a skew-symmetric matrix  $S = (S_{ik})$  such that the quadratic form  $S_{ik}x_iy_k$  is positive definite on L.

Let m be the dimension of L and let  $(X_{1\alpha}, \dots, X_{n\alpha}, Y_{1\alpha}, \dots, Y_{n\alpha})$   $(\alpha = 1, \dots, m)$  be a basis for L, that is, a set of m points  $(X_{\alpha}, Y_{\alpha})$  in L such that the coordinates of each point (x, y) in L are expressible uniquely in the form

(2) 
$$x_i = X_{i\alpha}u_{\alpha}, \qquad y_i = Y_{i\alpha}u_{\alpha} \qquad (\alpha = 1, \cdots, m).$$

Since the matrix  $(x_i, y_i)$  has rank 2 at each point  $(x, y) \neq (0, 0)$  on L, the matrix

 $\lambda X + \mu Y$ , where  $X = (X_{i\alpha})$  and  $Y = (Y_{i\alpha})$ , has rank m for every pair of real numbers  $\lambda$ ,  $\mu$ , not both zero. Conversely for every pair of matrices X, Y, having  $\lambda X + \mu Y$  of rank m when  $(\lambda, \mu) \neq (0, 0)$ , the corresponding linear space L defined by equations (2) is such that the matrix  $(x_i, y_i)$  has rank 2 at each point  $(x, y) \neq (0, 0)$  on L. Moreover by the use of equations (2) it is seen that  $S_{ik}x_iy_k = R_{\alpha\beta}u_\alpha u_\beta$ , where R = X'SY and X' is the transpose of X. Here and elsewhere it is understood that the symbol for an element of a matrix is obtained by adding a pair of subscripts to the symbol for the matrix. Theorem 3 is accordingly equivalent to the following theorem on matrices:

THEOREM 4. Let X, Y be two  $n \times m$ -dimensional matrices such that the matrix  $\lambda X + \mu Y$  has rank m for every pair of real numbers  $\lambda$ ,  $\mu$ , not both zero. There exists an n-rowed skew-symmetric matrix S such that the matrix R = X'SY is positive definite, that is, the quadratic form  $R_{\alpha\beta}u_{\alpha}u_{\beta}$   $(\alpha, \beta = 1, \dots, m)$  is positive definite. Here X' is the transpose of X. The matrix R in general will not be symmetric.

In order to prove Theorem 4 we first observe that in the proof of the theorem we may replace the matrices X, Y by  $X_1 = AXB$ ,  $Y_1 = AYB$ , where A and B are arbitrary nonsingular matrices of dimensions n and m respectively. For, the matrix  $\lambda X_1 + \mu Y_1$  has rank m when  $(\lambda, \mu) \neq (0, 0)$ . Moreover, suppose there exists a skew-symmetric matrix  $S_1$  such that the matrix  $X_1'S_1Y_1$  is positive definite. Let S be the skew-symmetric matrix  $A'S_1A$ , and set R = X'SY. Then the matrix  $B'RB = X_1'S_1Y_1$  is positive definite. It follows readily that R is positive definite and hence that the matrices X, Y can be replaced by  $X_1$ ,  $Y_1$  respectively. In fact X, Y can be replaced by matrices of the form

$$(3) X = \begin{pmatrix} I \\ O \end{pmatrix}, Y = \begin{pmatrix} C \\ D \end{pmatrix},$$

where I is the m-rowed identity matrix, C is an m-rowed matrix and O is the zero matrix. Suppose that X, Y are of this form. Then the condition that the matrix  $\lambda X + \mu Y$  be of rank m for all real numbers  $\lambda$ ,  $\mu$ , not both zero, is equivalent to the condition that the matrix

$$\begin{pmatrix} \lambda I - C \\ -D \end{pmatrix}$$

be of rank m for all real values of  $\lambda$ . Moreover the equation AXB = X holds for two nonsingular matrices A and B if and only if A is of the form

(5) 
$$A = \begin{pmatrix} B^{-1} & B^{-1}E \\ O & F \end{pmatrix}.$$

When A is of this form we have

(6) 
$$AYB = {B^{-1}(C + ED)B \choose FDB}.$$

In the sequel we shall assume that the matrices X, Y, A are of the forms (3) and (5). It will be understood that if n=m the matrices O, D, E, F do not appear in these matrices. Theorem 4 will be established by replacing Y by a matrix AYB having special properties. This will be done with the help of three lemmas, the first of which is the following:

LEMMA 1. Let  $\phi_1(\lambda)$ ,  $\cdots$ ,  $\phi_k(\lambda)$  be the elementary divisors of the matrix  $\lambda I - C$  and let  $\epsilon$  be an arbitrary constant different from zero. With an elementary divisor of the form  $\phi_i(\lambda) = [(\lambda - \alpha_i)^2 + \beta_i^2]^{m_i}$  ( $\beta_i \neq 0$ ) associate the  $2m_i$ -rowed matrix

(7) 
$$M_{i} = \begin{pmatrix} N_{i} & O & O & \cdots & O \\ \epsilon I & N_{i} & O & \cdots & O \\ O & \epsilon I & N_{i} & \cdots & O \\ \vdots & \vdots & \ddots & \ddots & \vdots \\ O & O & O & \cdots & \epsilon I & N_{i} \end{pmatrix}, \quad N_{i} = \begin{pmatrix} \alpha_{i} & -\beta_{i} \\ \beta_{i} & \alpha_{i} \end{pmatrix}, \quad I = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}.$$

With an elementary divisor of the form  $\phi_i(\lambda) = (\lambda - \alpha_i)^{m_i}$  associate the matrix

(8) 
$$M_{i} = \begin{pmatrix} \alpha_{i} & 0 & 0 & \cdots & 0 \\ 1 & \alpha_{i} & 0 & \cdots & 0 \\ 0 & 1 & \alpha_{i} & \cdots & 0 \\ \vdots & \vdots & \ddots & \ddots & \vdots \\ 0 & 0 & 0 & \cdots & 1 & \alpha_{i} \end{pmatrix}.$$

There exists a nonsingular m-rowed matrix B such that  $B^{-1}CB = C_1$ , where

(9) 
$$C_{1} = \begin{pmatrix} M_{1} & O & O & \cdots & O \\ O & M_{2} & O & \cdots & O \\ O & O & M_{3} & \cdots & O \\ \vdots & \vdots & \ddots & \ddots & \vdots \\ O & O & O & \cdots & M_{k} \end{pmatrix}.$$

For it is clear from the special forms (7) and (8) that  $\phi_i(\lambda) = |\lambda I - M_i|$ . If  $M_i$  is of the form (8), the  $(m_i - 1)$ -rowed minor obtained by deleting the first row and last column has determinant unity. Hence  $\phi_i(\lambda)$  is the only non-trivial invariant factor of  $\lambda I - M_i$  in this case. Suppose next that  $M_i$  is of the

form (7) and let  $\phi(\lambda)$  be the determinant of the minor of  $\lambda I - M_i$  obtained by deleting the first row and the last column. The values  $\phi(\alpha_i + i\beta_i), \phi(\alpha_i - i\beta_i)$  are different from zero since  $\epsilon \neq 0$ , as one readily verifies. It follows that  $\phi(\lambda)$  and  $\phi_i(\lambda)$  have no common factors and hence that  $\phi_i(\lambda)$  is the only non-trivial invariant factor of  $\lambda I - M_i$  when  $M_i$  is of the form (7). The matrices  $\lambda I - C_1$  and  $\lambda I - C$  accordingly have the same elementary divisors and hence the same invariant factors. The matrix  $C_1$  is therefore similar to C, that is,  $C_1$  is of the form  $B^{-1}CB$ , where B is nonsingular.  $\dagger$ 

Lemma 2. If the equation  $|\lambda I - C| = 0$  has no real roots, there exists a matrix  $C^*$  similar to C and a skew-symmetric matrix S such that the matrix  $SC^*$  is positive definite. The matrix  $SC^*$  is not in general symmetric.

For in this case the diagonal blocks in (9) are of the form (7). Consider  $C_1$  as a function  $C_1(\alpha, \beta, \epsilon)$  of the values  $\alpha_i$ ,  $\beta_i$ ,  $\epsilon$  described in Lemma 1. Let  $R(\epsilon) = C_1(0, \beta, \epsilon)$ , S = -R(0) and  $T = C_1(\alpha, 0, 0)$ . Then S and ST are skew-symmetric, S is nonsingular and  $C_1 = R(\epsilon) + T$ . Since ST is skew-symmetric one has  $u_{\alpha}S_{\alpha\beta}T_{\beta\gamma}u_{\gamma} = 0$   $(\alpha, \beta, \gamma = 1, \cdots, m)$  and

$$u_{\alpha}S_{\alpha\beta}(R_{\beta\gamma} + T_{\beta\gamma})u_{\gamma} = u_{\alpha}S_{\alpha\beta}R_{\beta\gamma}u_{\gamma}$$

for arbitrary values of (u). Since R(0) = -S and S is nonsingular, the last quadratic form is positive definite when  $\epsilon = 0$  and hence for a value  $\epsilon' \neq 0$ . The matrices S and  $C^* = R(\epsilon') + T$  have the properties described in the lemma.

LEMMA 3. Let C, D be the matrices appearing in the matrix (4) and suppose the matrix (4) has rank m for all real values of  $\lambda$ . There exists an  $m \times (n-m)$ -dimensional matrix E such that the equation

$$(10) |\lambda I - C - ED| = 0$$

has no real roots if the dimension m of C is even and a single real root if m is odd.

Let us begin by disposing of the case m=n. If m=n, the conclusion of Lemma 3 is that  $|\lambda I - C| \neq 0$  for all real  $\lambda$ . This is an immediate consequence of the hypothesis that matrix (4) has rank m=n for all real  $\lambda$ . Incidentally, if m=n, both m and n must be even; otherwise the equation  $|\lambda I - C| = 0$  would be of odd degree, and would therefore have a real root.

Suppose then that m < n. We shall consider first the case in which m = 2 and C, D are of the form

(11) 
$$C = \begin{pmatrix} a & 0 \\ c & b \end{pmatrix}, \qquad D = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} \text{ or } D = (d \ e).$$

<sup>†</sup> A. A. Albert, Modern Higher Algebra, University of Chicago Press, 1937, pp. 84-85.

If D is the two-rowed identity matrix, let

$$G = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}, \qquad E = G - C.$$

Then the equation (10) takes the form  $\lambda^2 + 1 = 0$  and has no real roots. Suppose next that D = (d e). We may suppose that b = 0 since this can be brought about by replacing  $\lambda$  by  $\lambda + b$ . Since the matrix (4) has rank 2 for  $\lambda = 0$  and  $\lambda = a$ , we have  $e \neq 0$  and  $ad + ec \neq 0$ . Choose numbers  $\alpha$ ,  $\beta$  such that

$$\alpha ce - \beta ae = -1, \quad \alpha d + \beta e + a = 0.$$

Then equation (10) with  $E' = (\alpha \beta)$  as the transpose of E reduces to  $\lambda^2 + 1 = 0$ . The lemma is accordingly true for the matrices (11).

To prove the lemma as stated, let h be the number of real roots of the equation  $|\lambda I - C| = 0$ , each root counted a number of times equal to its multiplicity. We may suppose that  $h \ge 2$ . It is sufficient to show that the matrix E can be chosen so that the equation (10) has exactly h-2 real roots. In view of equation (6) it is sufficient to prove this result when C, D are replaced respectively by matrices of the form  $C_1 = B^{-1}CB$ ,  $D_1 = FDB$ , where B and F are nonsingular. By virtue of Lemma 1 we may select B so that  $C_1$  is given by (9). In fact B may be chosen so that the matrices (8) corresponding to real roots of  $|\lambda I - C| = 0$  have higher subscripts than the matrices (7) corresponding to roots that are not real. It follows that, after a suitable choice of the matrix F in  $D_1$ , one has

$$C_1 = \begin{pmatrix} M & O \\ N & C_2 \end{pmatrix}, \qquad D_1 = \begin{pmatrix} U & O \\ V & D_2 \end{pmatrix},$$

where the matrices  $C_2$ ,  $D_2$  are the matrices C, D in (11). The matrix (4) with  $C = C_2$ ,  $D = D_2$  has rank 2, since otherwise the corresponding matrix with  $C = C_1$ ,  $D = D_1$  could not have rank m. It follows that there exists a matrix  $E_2$  such that the equation (10) with  $C = C_2$ ,  $D = D_2$ ,  $E = E_2$  has no real roots. Choosing

$$E = \begin{pmatrix} O & O \\ O & E_2 \end{pmatrix},$$

the equation (10) with  $C = C_1$ ,  $D = D_1$  reduces to the product

$$|\lambda I_{m-2}-M|\cdot|\lambda I_2-C_2-E_2D_2|=0,$$

where the subscript on I denotes its dimension. This equation has h-2 real roots and the proof of Lemma 3 is complete.

We are now in position to prove Theorem 4. As was seen above we can assume that the matrices X, Y are of the form (3). Since we can replace Y

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by the matrix (6) we can suppose, by virtue of Lemma 3, that the equation  $|\lambda I - C| = 0$  has no real roots if the dimension m of C is even and one real root if m is odd. Consider first the case in which m is even. Then by Lemma 2 there is a nonsingular matrix B and a skew-symmetric matrix S such that the matrix  $SC^*$  with  $C^* = B^{-1}CB$  is positive definite. Replace Y by AYB, where A is defined by equation (5) with E = 0 and F = I. Then C is replaced by  $C^*$  so that the matrix SC is now positive definite. Let  $S_1$  be the n-rowed skew-symmetric matrix

$$\begin{pmatrix} S & O \\ O & O \end{pmatrix}$$
.

The product  $X'S_1Y$  is then equal to SC and is accordingly positive definite. This proves Theorem 4 for the case in which m is even.

The case in which m is odd will be reduced to the case in which m is even. As was seen above we can assume that the equation  $|\beta I - C| = 0$  has only one real root  $\lambda = a$ . Choose nonsingular matrices B, F such that the matrices  $C_1 = B'CB$ ,  $D_1 = FDB$  are of the form

$$C_1 = \begin{pmatrix} M & O \\ O & a \end{pmatrix}, \qquad D_1 = \begin{pmatrix} U & b \\ V & O \end{pmatrix} \qquad (b \neq 0).$$

This choice is possible by virtue of Lemma 1 with  $M_k = (a)$  and the fact that the matrix (4) has rank m when  $\lambda = a$ . We can suppose that  $C = C_1$ ,  $D = D_1$ , since this result can be brought about by replacing Y by AYB, where A is given by (5) with E = O. Clearly m < n since D has at least one row. Let  $X_1, Y_1$  be the matrices

$$X_{1} = \begin{bmatrix} I & O & O \\ O & 1 & 0 \\ O & 0 & 1 \\ O & O & O \end{bmatrix}, \qquad Y_{1} = \begin{bmatrix} M & O & O \\ O & a & -b \\ U & b & a \\ V & O & O \end{bmatrix}$$

obtained by adding a suitable column to each of the matrices X and Y. Since the equations  $|\lambda I - M| = 0$ ,  $(\lambda - a)^2 + b^2 = 0$  have no real roots, the matrix  $\lambda X_1 - Y_1$  corresponding to (4) has rank m+1 for all real values of  $\lambda$ . As was seen in the last paragraph, there exists, since m+1 is even, a skew-symmetric matrix S such that the matrix  $X_1'SY_1$  is positive definite. The matrix  $X_2'SY_2$ , being a principal minor of  $X_1'SY_2$ , is also positive definite. This completes the proof of Theorem 4 and hence also of Theorem 3.

Incidentally, by an argument similar to that just made, it can be shown that if L is a maximal linear space having the properties described in Theorem 3 then its dimension is n if n is even and n-1 if n is odd.

3. Application to the calculus of variations. Consider the problem of minimizing the double integral

$$I = \int \int_A f(x, y, z_1, \cdots, z_n, p_1, \cdots, p_n, q_1, \cdots, q_n) dxdy,$$

where  $p_i = \partial z_i / \partial x$ ,  $q_i = \partial z_i / \partial y$ , in a class of subspaces

(12) 
$$z_i(x, y)$$
  $((x, y) \text{ on } A + B; i = 1, \dots, n)$ 

having a common boundary, where B is the boundary of A. The integrand f(x, y, z, p, q) is assumed to be continuous and to have continuous first and second partial derivatives on a region R of points (x, y, z, p, q). The subspaces (12) are assumed to be continuous, to have continuous second partial derivatives and to have their elements (x, y, z, p, q) in R. The boundary B of A is supposed to be a simply closed continuous curve having a piecewise continuously turning tangent. Weaker differentiability assumptions could be made.

LEMMA 4. Let  $S_{ik}(x, y) = -S_{ki}(x, y)$  (i,  $k = 1, \dots, n$ ) be arbitrary continuous functions having continuous first and second derivatives in a neighborhood of the set A + B. Then the integral

(13) 
$$J = \int \int_{A} g(x, y, z, p, q) dx dy$$

$$\equiv \int \int_{A} (1/2) \left\{ p_{i} \frac{\partial}{\partial y} (s_{ik} z_{k}) - q_{i} \frac{\partial}{\partial x} (s_{ik} z_{k}) \right\} dx dy$$

has the same value for all subspaces (12) here considered.

For by virtue of Green's theorem the value of 2J is given by the formula

$$2J = \int\!\!\int_{A} \left\{ \frac{\partial}{\partial y} \left( S_{ik} p_{i} z_{k} \right) - \frac{\partial}{\partial x} \left( S_{ik} q_{i} z_{k} \right) \right\} dx dy = \int_{B} z_{k} S_{ki} dz_{i},$$

and hence is completely determined by the common boundary of the subspaces (12).

Let E be a minimizing subspace. Then E must satisfy the condition of Legendre,\* that is, at each element (x, y, z, p, q) on E the inequality

(14) 
$$P(\xi,\eta) = f_{\nu_i \nu_k} \xi_i \xi_k + 2 f_{\nu_i q_k} \xi_i \eta_k + f_{q_i q_k} \eta_i \eta_k \ge 0$$

must hold for every set  $(\xi, \eta)$  whose  $n \times 2$ -dimensional matrix has rank 1. If

<sup>\*</sup> Graves, The Weierstrass condition for multiple integral variation problems, Duke Mathematical Journal, vol. 5 (1939), pp. 656-660. Graves proves only the Weierstrass condition. The condition of Legendre is a well known consequence of that of Weierstrass.

the condition (14) with the equality excluded holds on E for every set  $(\xi, \eta)$  whose matrix  $(\xi_i, \eta_i)$  has rank 1, then E will be said to satisfy the *strengthened* condition of Legendre. It should be emphasized that in this strengthened condition of Legendre we require that the inequality  $P(\xi, \eta) > 0$  hold only for sets  $(\xi, \eta)$  whose  $n \times 2$ -dimensional matrix  $(\xi_i, \eta_i)$  has rank 1 and not for all sets  $(\xi, \eta) \neq (0, 0)$ . We have the following:

THEOREM 5. Let E be a subspace (12) satisfying the strengthened condition of Legendre. There exists an invariant integral of the form (13) such that for the function F = f + g the inequality

$$(15) F_{p_i p_i} \xi_i \xi_k + 2F_{p_i q_i} \xi_i \eta_k + F_{q_i q_i} \eta_i \eta_k > 0$$

holds at each element (x, y, z, p, q) on E for every set  $(\xi, \eta) \neq (0, 0)$ .

For by Theorem 2 there exists for each point on E a skew-symmetric matrix  $(S_{ik})$  such that the quadratic form  $P(\xi, \eta) + S_{ik}\xi_i\eta_k$  is positive definite at this point and hence in a neighborhood of this point. It follows readily that the set A+B can be covered by a finite number of circles  $C_1, \dots, C_m$ having a common radius r such that to the circle  $C_{\alpha}'$  of radius 2r concentric with the circle  $C_{\alpha}$  there is associated a skew-symmetric matrix  $(S_{ik}^{\alpha})$  of constant elements such that the form  $P_{\alpha}(\xi, \eta) = P(\xi, \eta) + S_{ik}^{\alpha} \xi_i \eta_k$  is positive definite at each point on E whose projection in the xy-plane is in  $C_{\alpha}'$ . We propose now to construct a set of functions  $S_{ik}(x, y) = -S_{ki}(x, y)$  having continuous second derivatives on a neighborhood of A+B and such that the form  $P(\xi, \eta) + S_{ik}(x, y)\xi_i\eta_k$  is positive definite on E. To this end let h(t) be a function having continuous second derivatives and such that h(t) = 0 when  $t \le r$ , h(t) = 1 when  $t \ge 2r$ , and 0 < h(t) < 1 when r < t < 2r. Let  $(x_{\alpha}, y_{\alpha})$  be the center of the circle  $C_{\alpha}$ , and let  $d_{\alpha}(x, y)$  be the distance between the points (x, y)and  $(x_{\alpha}, y_{\alpha})$ . Set  $h_{\alpha}(x, y) = h[d_{\alpha}(x, y)]$ . It should be observed that the product  $h_1h_2 \cdots h_m$  is identically zero on the set A+B. This follows because  $h_\alpha \equiv 0$ on  $C_{\alpha}$  and the circles  $C_1, \dots, C_m$  cover A+B. Denote by  $S_{ik}(x, y)$  the function

$$(1-h_1)S_{ik}^1+h_1(1-h_2)S_{ik}^2+\cdots+h_1h_2\cdots h_{m-1}(1-h_m)S_{ik}^m$$

where  $S_{ik}^{\alpha}$   $(\alpha = 1, \dots, m)$  are the constants described above. The quadratic form  $P(\xi, \eta) + S_{ik}(x, y)\xi_i\eta_k$  determined by these functions is easily seen to be identical with the sum

$$(1-h_1)P_1(\xi,\eta)+h_1(1-h_2)P_2(\xi,\eta)+\cdots+h_1h_2\cdots h_{m-1}(1-h_m)P_m(\xi,\eta),$$

where  $P_{\alpha}(\xi, \eta)$  is the quadratic form associated with the circle C'. Each term in this sum is positive or zero, since for each integer  $\alpha$  the product

 $(1-h_{\alpha})P_{\alpha}(\xi, \eta)$  is positive definite at points on E corresponding to points in  $C'_2$  and identically zero elsewhere. Moreover the terms do not vanish simultaneously when  $(\xi, \eta) \neq (0, 0)$ . The quadratic form  $P(\xi, \eta) + S_{ik}\xi_i\eta_k$  is accordingly positive definite on E.

We now use these functions  $S_{ik}$  to define the function g of equation (13). From the definition of g we at once obtain

$$g(x, y, z, p, q) = (1/2) \left[ (\partial S_{ik}/\partial y) p_i z_k + (\partial S_{ki}/\partial x) q_i z_k \right] + S_{ik} p_i q_k.$$

This shows that the Legendre quadratic form (15) for g is  $S_{ik}\xi_i\eta_k$ . Hence if F=f+g the Legendre form (15) for F is exactly the positive definite form  $P(\xi, \eta)+S_{ik}\xi_i\eta_k$ . This establishes Theorem 5.

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