REGULAR BOUNDARY CONDITIONS FOR ORDINARY DIFFERENTIAL OPERATORS

BY STEPHEN SALAFF

Introduction. The Birkhoff theory for linear differential operators of order m treats operators determined by a nonselfadjoint differential expression and two-point boundary conditions. In his classic of 1908 [1], Birkhoff gave determinantal inequalities, involving the leading coefficients in the boundary forms, from which asymptotic eigenfunction estimates and an expansion theorem follow. Stone [14] then showed that the Birkhoff expansion is in a sense equivalent to the Fourier expansion.

In the present paper we focus on a question of Kamke [5, p. 197], and show that if m is even, selfadjoint boundary conditions meet the determinantal inequalities (which are somewhat simplified in the process). This result is not unexpected, for if the differential expression is formally selfadjoint, the operator determined is selfadjoint and an expansion theorem (such as the Fourier expansion) is available.

The difficulties in obtaining the result are (linear) algebraic ones. Consider the matrix $(f(x_j+x_k))_{j,k=1}^n$, where $x_1 < x_2 < \cdots < x_n$. Let f satisfy the functional equation $f(x+y)=g(x)g(y)[h(x)+h(y)]^{-1}$, where g and h are suitably restricted. In Lemma III we use the function $f(x)=\csc x$ to show that the given matrix, and a class of others related to its compounds, are positive definite. This development is central to the proof of our theorem.

I. Preliminaries.

1. Notation. In this section C^{2m} denotes the 2m dimensional complex vector space, with inner product $(x, y) = \sum_{j=1}^{2m} x_j \bar{y}_j$, $x = (x_1, \dots, x_{2m})$, $y = (y_1, \dots, y_{2m})$. The Hilbert space $H = L^2(0, 1)$ is the space of Lebesgue square-integrable functions in the closed unit interval, with inner product

$$(u, v) = \int_0^1 u(t) \overline{v(t)} dt.$$

The subset H_1 of H consists of those functions having absolutely continuous derivatives to order m-1, with mth derivative in H. For $u \in H_1$ we define $\hat{u} \in C^{2m}$ by $\hat{u} = (u(0), u'(0), \dots, u^{(m-1)}(0), u(1), \dots, u^{(m-1)}(1))$, where $u^{(j)}$ is the jth derivative of u.

The transposed complex conjugate of the matrix B is written B^* . The large cross X which occurs in matrices and equations of this chapter will indicate the presence of certain terms whose value is not of present concern.

2. The differential operator L. The linear differential expression l, of arbitrary order m, is defined in H_1 by

$$lu = \sum_{j=0}^m p_{m-j} u^{(j)},$$

where $p_0 \equiv 1$ and each p_j , together with its derivatives to order m-j, is in H. The formal adjoint expression l^* is given for $v \in H_1$ by

$$l^*v = \sum_{j=0}^m (-1)^j (\bar{p}_{m-j}v)^{(j)}.$$

Lagrange's formula [2, p. 285] is

$$(lu, v) - (u, l^*v) = (A\hat{u}, \hat{v}).$$

Here A is a nonsingular $2m \times 2m$ matrix of the form

There will be m linearly independent, homogeneous, two-point boundary conditions

$$bu \equiv \sum_{j=0}^{m-1} B_j^0 u^{(j)}(0) + B_j^1 u^{(j)}(1) = 0.$$

The complex column vectors B_i^0 , B_i^1 are of length m. Independence requires that rank B=m, where B is the $m \times 2m$ matrix

$$(B_0^0,\ldots,B_{m-1}^0,B_0^1,\ldots,B_{m-1}^1).$$

Thus bu=0 can also be written $B\hat{u}=0$. We shall refer to the left side of a boundary condition as a boundary form. Such a form is said to be of order k, $0 \le k \le m-1$, if the highest order derivative of u appearing in it (at either endpoint 0 or 1) is the kth derivative.

Any boundary conditions obtained from $B\hat{u}=0$ by taking linear combinations (left-multiplying B by a nonsingular $m \times m$ matrix T) are termed 'equivalent' conditions, since $B\hat{u}=0$ if and only if $TB\hat{u}=0$. We choose a particular T, corresponding to the 'normalized' boundary forms, which will now be described.

Let $r_{m-1} = \text{rank } (B_{m-1}^0, B_{m-1}^1)$. Then r_{m-1} is the number (0, 1 or 2) of independent boundary forms of order m-1. There is a nonsingular $m \times m$ matrix T_{m-1} , and there are two short column vectors b_{m-1}^0, b_{m-1}^1 , of length r_{m-1} , such that

$$T_{m-1}(B_{m-1}^0, B_{m-1}^1) = \begin{bmatrix} 0 & 0 \\ \vdots & \vdots \\ 0 & 0 \\ b_{m-1}^0 & b_{m-1}^1 \end{bmatrix}.$$

If $r_{m-1}=0$, this matrix consists of zeros only, T_{m-1} can be taken as the identity, and the columns b_{m-1}^0 , b_{m-1}^1 will be empty. Now cross off the last r_{m-1} rows of $T_{m-1}(B_{m-2}^0, B_{m-2}^1)$, and let r_{m-2} be the rank of the matrix remaining. Again, r_{m-2} is 0, 1 or 2, and there are T_{m-2} , b_{m-2}^0 , b_{m-2}^1 such that

$$T_{m-2}T_{m-1}(B_{m-2}^{0}, B_{m-2}^{1}) = \begin{bmatrix} 0 & 0 \\ \vdots & \vdots \\ 0 & 0 \\ b_{m-2}^{0} & b_{m-2}^{1} \\ X & X \end{bmatrix}.$$

The crosses indicate the presence of r_{m-1} rows.

Continuing in this way till no further forms remain, we obtain the normalized boundary forms, with $T = T_0 \cdots T_{m-2} T_{m-1}$. The normalized boundary matrix, hereafter referred to as B, will have the form

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$$B$$
, will have the form
$$\begin{bmatrix} b_0^0 & 0 & 0 & \cdots & 0 & b_0^1 & 0 & 0 & \cdots & 0 \\ b_1^0 & 0 & \cdots & 0 & & b_0^1 & 0 & \cdots & 0 \\ & & \ddots & & & & & \ddots & \vdots \\ X & & 0 & & & X & & 0 \\ & & & b_{m-1}^0 & & & & b_{m-1}^1 \end{bmatrix}$$

with rank $(b_j^0, b_j^1) = r_j$ and $\sum_{j=0}^{m-1} r_j = m$. The normalized conditions read $b_i^0 u^{(j)}(0) + \mathbf{X} + b_j^1 u^{(j)}(1) + \mathbf{X} = 0$, j = 0, 1, ..., m-1.

The crosses take the place of lower order terms at 0 and at 1.

In illustration, if m=4, $r_0=r_3=1$, $r_1=2$, $r_2=0$, then

We now define the differential operator L, with domain $D_L = \{u \in H_1 : bu = 0\}$, by setting Lu = lu. Then L is a densely defined operator in H, and is closed [4].

3. The adjoint, and selfadjointness conditions. Associated with each (normalized) boundary matrix B is its m dimensional null space S, spanned by m linearly independent solutions of Bx=0. Conversely, given an m dimensional subspace of C^{2m} , there is an $m \times 2m$ matrix B, unique to within left multiplication by an arbitrary nonsingular matrix T, which annihilates this subspace. A one to one correspondence between classes of equivalent boundary conditions and the subspaces associated with their matrices is established in this way. The matrix A of the previous section maps S onto another m dimensional subspace AS. Let B' be a boundary matrix corresponding to $(AS)^{\perp}$, the orthogonal complement of AS in C^{2m} . We define a differential operator L' with $D_{L'} = \{v \in H_1 : B'\hat{v} = 0\}$ by putting $L'v = l^*v$. Then $L' = L^*$, the operator adjoint to L [4].

When $D_L = D_{L^*}$, i.e. when $AS = S^{\perp}$, we say that "The conditions bu = 0 are selfadjoint with respect to l" and abbreviate this statement by "S.A." We also speak of such boundary conditions as selfadjoint, or of selfadjoint type, with the dependence on l being understood.

LEMMA 1. S.A. if and only if $BA^{-1}B^*=0$.

Proof. Since S is the null space of B we have $S^{\perp} = R(B^*)$, the range of B^* , and $AS = S^{\perp}$ if and only if $S = R(A^{-1}B^*)$, if and only if $BA^{-1}B^* = 0$.

Referring to the forms of A and B in 2, we see that A^{-1} , B^* and $BA^{-1}B^*$ are given, in this order, by the following three matrices.

$$BA^{-1}B^* =$$

In the example,

$$B^* = \begin{bmatrix} \bar{b}_0^0 & \times & \times & \times \\ 0 & (b_1^{0*}) & \times \\ 0 & 0 & 0 & \times \\ 0 & 0 & 0 & \bar{b}_3^0 \\ \bar{b}_0^1 & \times & \times & \times \\ 0 & (b_1^{1*}) & \times \\ 0 & 0 & 0 & \times \\ 0 & 0 & 0 & \bar{b}_3^1 \end{bmatrix},$$

$$BA^{-1}B^* = \begin{bmatrix} 0 & 0 & 0 & -b_0^0 \bar{b}_3^0 + b_0^1 \bar{b}_3^1 \\ 0 & 0 & 0 & \times \\ 0 & 0 & 0 & \times \\ 0 & 0 & 0 & \times \\ b_3^0 \bar{b}_0^0 - b_3^1 \bar{b}_0^1 & \times & \times \end{bmatrix}.$$
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The particular consequence of S.A. we use is

(1)
$$b_j^0 b_{m-1-j}^{0*} = b_j^1 b_{m-1-j}^{1*}, \quad j = 0, 1, ..., m-1.$$

It will be understood (as the example indicates) that the jth of these equations is absent if $r_i = 0$. When $r_j = 1$ we interpret b_j^0 , b_j^1 as scalars.

The equations (1) are the full extent of the information we would obtain if $lu=u^{(m)}$, for in this case all terms symbolized by crosses in A and $BA^{-1}B^*$ vanish.

LEMMA 2. If S.A., then for j = 0, 1, ..., m-1,

- (i) if $r_j = r_{m-1-j} = 1$, then $b_j^0 \bar{b}_{m-1-j}^0 = b_j^1 \bar{b}_{m-1-j}^1$,
- (ii) if $r_i = 2$ then $r_{m-1-i} = 0$,
- (iii) the cases $r_j = r_{m-1-j} = 0$, and $r_j = 0$, $r_{m-1-j} = 1$ cannot occur.

Proof. Equation (1) implies (i).

If $r_i = 2$ there is a nonsingular 2×2 matrix C such that

$$Cb_j^0 = \begin{pmatrix} 1 \\ 0 \end{pmatrix}$$
 and $Cb_j^1 = \begin{pmatrix} 0 \\ 1 \end{pmatrix}$.

Then from (1),

$$\begin{pmatrix} 1 \\ 0 \end{pmatrix} b_{m-1-j}^{0*} = \begin{pmatrix} 0 \\ 1 \end{pmatrix} b_{m-1-j}^{1*}.$$

Hence (ii), and the implication that $r_j + r_{m-1-j} \le 2$ for each j. Now since $\sum_{j=0}^{m-1} r_j = m$, we must have $r_j + r_{m-1-j} = 2$, and the cases in (iii) are excluded.

For use in §II, where S.A. is assumed, we suppose that k+2h=m and that $r_{f_0} = \cdots = r_{f_{h-1}} = 2$, $r_{g_0} = \cdots = r_{g_{k-1}} = 1$. Let the sequences $\{f_j\}$, $\{g_j\}$ be strictly increasing. Then of course $g_j + g_{k-1-j} = m-1$, $j = 0, \ldots, k-1$.

Also, without loss of generality the normalization can be accomplished so that

$$(b_j^1, b_j^0) = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \quad j = f_0, \ldots, f_{h-1}.$$

4. The Birkhoff expansion theorem. Birkhoff [1] defined regular boundary conditions. These entail the nonvanishing of two $m \times m$ determinants, to be defined in §II, depending only on b_j^0 , b_j^1 but not on the expression l. We write "R" for "The boundary conditions bu=0 are regular." Under assumption R, it can be shown [3], [10] that (i) The eigenvalues of L form a denumerable sequence $\{\lambda_j\}$ with no finite accumulation point. The adjoint boundary conditions are also regular [8], and the eigenvalues of L^* are $\{\bar{\lambda}_j\}$. (ii) Explicit asymptotic estimates, for large $|\lambda|$, can be obtained for the eigenvalues and corresponding eigenfunctions. (iii) Under the additional assumption that each eigenvalue is simple, any function f in D_L can be expanded in a uniformly convergent series

$$f = \sum_{i=1}^{\infty} (f, z_i) y_i.$$

The y_j and z_j are the eigenfunctions of L and L^* corresponding to the eigenvalues λ_j , λ_j , and $(y_j, z_j) = 1$. (iv) If the boundary conditions are 'strongly regular' (another determinantal condition if m is even, one implied by the additional assumption of (iii)), the sequence $\{y_j\}$ is a Riesz basis in H, and we can normalize so that $(y_j, z_k) = \delta_{ik}$ [7], [2, p. 311].

5. Statement of our theorem; remarks; acknowledgement. Writers on the Birkhoff expansion theorem have assumed R when necessary for the analysis. The nature of regular boundary conditions remains largely unexplored. In particular, Kamke [5, p. 197] noted that the question of whether selfadjoint boundary conditions are regular appeared open, except for m=2 where the answer is 'yes'. In this paper we prove

THEOREM. If m is even, and if the boundary conditions are selfadjoint with respect to the differential expression $u^{(m)}$ then they are regular for the expression l.

By virtue of this theorem, and the remark after equations (1), we can assert that if m is even then S.A. implies R.

§II contains a simplification of the regularity determinants, and exploitation of the selfadjointness hypothesis to cast them into one standard form. This is done for both even and odd m. In §III we construct a class of positive definite matrices which will be central to the proof in §IV that the standard form is nonvanishing for even m. Before proceeding, we make two remarks.

- (a) A proof of the theorem for odd m has not yet been found, except for m=3 and 5. We are hopeful, however, that the present methods can be extended to this case also (1).
- (b) The theorem shows that if m is even and L is normal then R. This follows from the fact [9, p. 33] that $D_L = D_{L^{\bullet}}$ for (closed) normal operators L. Conditions for normality are listed in [12].

I would like to thank Professor H. O. Cordes for suggesting the problem under consideration, and for showing me the close relation between Lemma III and the proof of the theorem. I am also grateful to Professors Donald Sarason and David Spring for many useful discussions.

II. Simplification of the regularity determinants.

1. Replacement of ω_j by ε_j . Defining regularity [3], we must distinguish case (E), m even, and case (O), m odd. In (E), m=2q and the boundary conditions bu=0 are regular if $\theta_0 \neq 0$ and $\theta_2 \neq 0$ in the identity

In (O), m=2q+1 and we have regularity if $\theta_0 \neq 0$, $\theta_1 \neq 0$ in the identity

$$\theta_{0} + \theta_{1}s = \begin{bmatrix} \omega_{0}^{0}b_{0}^{1} & \cdots & \omega_{q-1}^{0}b_{0}^{1} & \omega_{q}^{0}(b_{0}^{0} + b_{0}^{1}s) & \omega_{q+1}^{0}b_{0}^{0} & \cdots & \omega_{m-1}^{0}b_{0}^{0} \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ \omega_{0}^{m-1}b_{m-1}^{1} & \cdots & \omega_{q-1}^{m-1}b_{m-1}^{1} & \omega_{q}^{m-1}(b_{m-1}^{0} + b_{m-1}^{1}s) & \omega_{q+1}^{m-1}b_{m-1}^{0} & \cdots & \omega_{m-1}^{m-1}b_{m-1}^{0} \end{bmatrix}.$$

The sequence of distinct unit roots $\omega_0, \omega_1, \ldots, \omega_{m-1}$ is a permutation of the naturally ordered sequence $\varepsilon_0, \varepsilon_1, \ldots, \varepsilon_{m-1}$, with $\varepsilon_j = e^{2\pi i j/m}$: $\omega_0 = \varepsilon_0 = 1$, $\omega_1 = \varepsilon_{m-1}$, $\omega_2 = \varepsilon_1, \omega_3 = \varepsilon_{m-2}, \ldots, \omega_{2j} = \varepsilon_j, \omega_{2j+1} = \varepsilon_{m-1-j}, \ldots$

We are going to show that the numbers ω_j can be replaced by ε_j in the determinants θ_0 , θ_1 and θ_2 . In (E) this will enable us to prove that $\theta_0 = \pm \theta_2$. A simplification will also be possible in (O). We define the standard form

$$\theta = \begin{vmatrix} \varepsilon_0^0 b_0^1 & \cdots & \varepsilon_{q-1}^0 b_0^1 & & \varepsilon_q^0 b_0^0 & \cdots & \varepsilon_{m-1}^0 b_0^0 \\ \vdots & & \vdots & & \vdots & & \vdots \\ \varepsilon_0^{m-1} b_{m-1}^1 & \cdots & \varepsilon_{q-1}^{m-1} b_{m-1}^1 & & \varepsilon_q^{m-1} b_{m-1}^0 & \cdots & \varepsilon_{m-1}^{m-1} b_{m-1}^0 \end{vmatrix}.$$

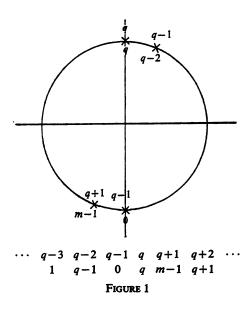
This definition is operative for both (E) and (O). The first q columns of θ involve b^1 (by which is meant b_j^1 , $j=0, 1, \ldots, m-1$). In (E), the last q involve b^0 and in (O) this number is q+1. Write $\theta=\theta(b^0,b^1)$ when it is necessary to emphasize that the first argument appears on the right, the second on the left.

LEMMA. R if and only if (E) $\theta \neq 0$; (O) $\theta(b^0, b^1) \neq 0$ and $\theta(b^1, b^0) \neq 0$.

⁽¹⁾ Added in proof. We have recently obtained a proof for odd m.

Proof. We begin with (E). That fact that b^1 appears only on the left and b^0 only on the right of θ_0 will not be altered if we execute certain interchanges among the first q columns, and also among the last q. The particulars depend on whether q is even or odd.

(a) q is even. The interchange is indicated in the following scheme. The first line consists of the column subscripts of the ω_j prior to the interchange, while the second contains the new subscripts (Figure 1).



(b) q odd (Figure 2).

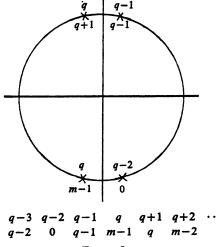


FIGURE 2

The new ordering, as shown by the integers outside the circumference in Figures 1 and 2, is cyclic in the counterclockwise sense. Putting

$$\phi = \omega_{q-1}$$
, q even,
= ω_{q-2} , q odd,

we see that θ is obtained from θ_0 by interchanging columns and replacing each ε_j by $\phi \varepsilon_j$. Hence

$$\theta_0 = \pm \theta.$$

In θ_2 , first interchange the center columns q-1 and q, so that b^1 appears only on the left, and b^0 only on the right.

(a) q even

$$0 \quad q-2 \quad q-1 \quad q \quad q+1 \quad q+2 \quad \cdots$$
 $0 \quad q-2 \quad m-1 \quad q-1 \quad m-2 \quad q$

(b) q odd

$$\cdots$$
 $q-3$ $q-2$ $q-1$ q $q+1$ $q+2$ \cdots $q-1$ 1 q 0 $q+1$ $m-1$

and again

$$\theta_2 = \pm \theta.$$

Taking up (O), we look first at θ_0 .

(a) q even

$$\cdots q-3 q-2 q-1 q q+1 q+2 \cdots 1 q-1 0 q m-1 q+1$$

(b) q odd

Passing finally to θ_1 , we transfer the last q columns, those containing b^0 , to the left.

(a) q even

(3)

$$\cdots q-3 q-2 q-1 q q+1 q+2 \cdots q+1 m-2 q m-1 q-1 0$$

(b) q odd

$$m-2$$
 $q-1$ q $q+1$ $q+2$ \cdots $q-2$ $q+1$ $m-1$ q q $q+1$

(4)
$$\theta_1 = \pm \theta(b^1, b^0).$$

Equations (1)-(4) establish the lemma.

2. Treatment of $r_i \neq 1$. We must come to grips with the fact that some of the matrices (b_j^1, b_j^0) may be empty. Consult §I.3 for the definition of f_j and g_j . Let L be the ordered set $\{g_0, \ldots, g_{l-1}\}$, where k=2l in (E) and k=2l+1 in (O). Also, let

$$(b_{g,i}^1, b_{g,i}^0) = (\beta_{g,i}, \alpha_{g,i}), |\alpha_{g,i}| + |\beta_{g,i}| > 0, \quad j = 0,1,\ldots,k-1.$$

Shift all the rows of θ corresponding to the f_i to the top, so that

all the rows of
$$\theta$$
 corresponding to the f_j to the top, so that
$$\theta = \pm \begin{bmatrix} \varepsilon_{f_0}^0 & \cdots & \varepsilon_{f_0}^{q-1} & 0 & \cdots & 0 \\ 0 & \cdots & 0 & e_{f_0}^q & \cdots & \varepsilon_{f_0}^{m-1} \\ & \cdots & & & & & & & & & & \\ \varepsilon_{f_{h-1}}^0 & \cdots & \varepsilon_{f_{h-1}}^q & 0 & \cdots & 0 \\ 0 & \cdots & 0 & \varepsilon_{f_{h-1}}^q & \cdots & \varepsilon_{f_{h-1}}^{m-1} \\ \beta_{g_0} \varepsilon_{g_0}^0 & \cdots & \beta_{g_0} \varepsilon_{g_0}^{q-1} & \alpha_{g_0} \varepsilon_{g_0}^q & \cdots & \alpha_{g_0} \varepsilon_{g_0}^{m-1} \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ \beta_{g_{k-1}} \varepsilon_{g_{k-1}}^0 & \cdots & \beta_{g_{k-1}} \varepsilon_{g_{k-1}}^{q-1} & \alpha_{g_{k-1}} \varepsilon_{g_{k-1}}^q & \cdots & \alpha_{g_{k-1}} \varepsilon_{g_{k-1}}^{m-1} \end{bmatrix}$$

LEMMA.

$$\theta = c \left[\prod_{0 \le i < i' \le h-1} (\varepsilon_{f_{j'}} - \varepsilon_{f_j})^2 \prod_{i=0: i'=0}^{k-1; h-1} (\varepsilon_{g_{j'}} - \varepsilon_{f_j}) \right] \theta_L,$$

where |c|=1 and

$$\theta_L = \begin{vmatrix} \beta_{g_0} \varepsilon_{g_0}^0 & \cdots & \beta_{g_0} \varepsilon_{g_0}^{l-1} & \alpha_{g_0} \varepsilon_{g_0}^q & \cdots & \alpha_{g_0} \varepsilon_{g_0}^{m+l-q-1} \\ \vdots & \vdots & \vdots & & \vdots \\ \beta_{g_{k-1}} \varepsilon_{g_{k-1}}^0 & \cdots & \beta_{g_{k-1}} \varepsilon_{g_{k-1}}^{l-1} & \alpha_{g_{k-1}} \varepsilon_{g_{k-1}}^q & \cdots & \alpha_{g_{k-1}} \varepsilon_{g_{k-1}}^{m+l-q-1} \end{vmatrix}.$$

Proof. In the formula for θ , empty products, occurring if h=0 or 1, or if k=0, will be interpreted as equal to one. We will refer to the columns of θ by the exponent of ε_{g_t} they carry. Perform the following column operations on θ . Multiply column q-1-j by ε_{f_0} and subtract from column q-j, for $j=1,\ldots,q-1$. Multiply column m-1-j by ε_{f_0} and subtract from column m-j, for $j=1,\ldots,m-1-q$.

Expand the resulting determinant from its first two rows, and then extract the factor

$$\prod_{j=1}^{k-1} (\varepsilon_{f_j} - \varepsilon_{f_0})^2 \prod_{j=0}^{k-1} (\varepsilon_{g_j} - \varepsilon_{f_0}).$$

The determinant remaining is of order m-2. We got it from θ by crossing out rows 0 and 1, and columns q-1 and m-1. This process is to be carried out h times, leading to the lemma. Thus θ_L is obtained from θ by crossing out the first 2h rows, and columns $q-j, m-j, j=1, \ldots, h$. We write $\theta_L = \theta_L(\alpha, \beta)$ to stress that the columns on the right contain the α_{g_i} .

COROLLARY. The boundary conditions bu=0 are regular if and only if

- (E) $\theta_L \neq 0$;
- (O) $\theta_L(\alpha, \beta) \neq 0$ and $\theta_L(\beta, \alpha) \neq 0$.

The corollary follows from the above and Lemma 1.

3. Use of selfadjoint type. If S.A. then

(5)
$$\alpha_{g_j}\bar{\alpha}_{g_{k-1-j}} = \beta_{g_j}\bar{\beta}_{g_{k-1-j}}, \quad j = 0, 1, \ldots, k-1,$$

an immediate consequence of Lemma I.2.

Once more, we are forced to distinguish cases (E) and (O).

In (E), k=2l and

$$\alpha_{g_i}\bar{\alpha}_{g_{2l-1-i}} = \beta_{g_i}\bar{\beta}_{2l-1-i}, \quad j = 0, 1, ..., l-1.$$

Multiply the last row of θ_L by $\bar{\alpha}_{g_0}$ and then remove the factor $\beta_{g_{2l-1}}$ from that row. Next multiply the same row by $\bar{\beta}_{g_0}$ and factor $\alpha_{g_{2l-1}}$. Writing θ'_L for the determinant whose last row is

$$\bar{\alpha}_{g_0}\varepsilon_{m-1-g_0}^0\cdot\cdot\cdot\bar{\alpha}_{g_0}\varepsilon_{m-1-g_0}^{l-1}\quad \bar{\beta}_{g_0}\varepsilon_{m-1-g_0}^{q}\cdot\cdot\bar{\beta}_{g_0}\varepsilon_{m-1-g_0}^{q+l-1}$$

and which agrees with θ_L in all other rows, we have

$$\bar{\alpha}_{g_0}\theta_L = \beta_{g_{2l-1}}\theta_L'$$
 and $\bar{\beta}_{g_0}\theta_L = \alpha_{g_{2l-1}}\theta_L'$

Thus

$$(|\alpha_{g_0}| + |\beta_{g_0}|) |\theta_L| = (|\alpha_{g_{2l-1}}| + |\beta_{g_{2l-1}}|) |\theta_L'|.$$

Doing this to each of the bottom l rows of θ_L , and abbreviating

$$\delta_L = \prod_{j=0}^{l-1} (|\alpha_{g_{2l-1-j}}| + |\beta_{g_{2l-1-j}}|)/(|\alpha_{g_j}| + |\beta_{g_j}|),$$

so that $\delta_L > 0$, we get

$$|\theta_L| = \delta_L |D|,$$

where

$$D = \begin{pmatrix} \beta_{g_0} \varepsilon_0^{g_0} & \cdots & \beta_{g_0} \varepsilon_{l-1}^{g_0} & \alpha_{g_0} \varepsilon_q^{g_0} & \cdots & \alpha_{g_0} \varepsilon_{q+l-1}^{g_0} \\ \vdots & & \vdots & & \vdots & & \vdots \\ \beta_{g_{l-1}} \varepsilon_0^{g_{l-1}} & \cdots & \beta_{g_{l-1}} \varepsilon_{l-1}^{g_{l-1}} & \alpha_{g_{l}} \varepsilon_q^{g_0} & \cdots & \alpha_{g_0} \varepsilon_{q+l-1}^{g_0} \\ \vdots & & & \vdots & & \vdots \\ \bar{\alpha}_{g_{l-1}} \varepsilon_0^{g_{l-1}} & \cdots & \bar{\alpha}_{g_{l-1}} \varepsilon_{l-1}^{g_{l-1}} & \alpha_{g_{l-1}} \varepsilon_q^{g_1} & \cdots & \alpha_{g_{l-1}} \varepsilon_{q+l-1}^{g_1} \\ \vdots & & & \vdots & & \vdots \\ \bar{\alpha}_{g_0} \varepsilon_0^{2q-1-g_0} & \cdots & \bar{\alpha}_{g_0} \varepsilon_{l-1}^{2q-1-g_0} & \bar{\beta}_{g_0} \varepsilon_{q+l-1}^{2q-1-g_0} & \cdots & \bar{\beta}_{g_0} \varepsilon_{q+l-1}^{2q-1-g_0} \end{pmatrix}$$

Combining this with Corollary 2, we have

PROPOSITION E. Boundary conditions of selfadjoint type are regular in case (E) if for arbitrary complex numbers α_j , β_j , subject only to the restriction that $|\alpha_j| + |\beta_j| > 0$, $j \in L$, we have $D \neq 0$.

In (O), m=2q+1, we have (5), and in addition

$$\alpha_a\bar{\alpha}_a=\beta_a\bar{\beta}_a\neq 0.$$

It will be convenient to assume $|\alpha_q| = |\beta_q| = 1$, and put $\alpha_q = e^{i\psi}$, $\beta_q = e^{i\chi}$, ψ and χ real. Employing the method of (E), and the δ_L defined there,

$$|\theta_L(\alpha, \beta)| = \delta_L |D'(\alpha, \beta)|,$$

where

$$D'(\alpha,\beta) = \begin{vmatrix} \beta_{g_0} \varepsilon_0^{g_0} & \cdots \\ \vdots & \vdots \\ \beta_{g_{l-1}} \varepsilon_0^{g_{l-1}} & \cdots \\ e^{i\psi} \varepsilon_0^{g} & \cdots & e^{i\psi} \varepsilon_{l-1}^{q} & e^{i\chi} \varepsilon_q^{q} & \cdots & e^{i\chi} \varepsilon_{q+l}^{q} \\ \bar{\alpha}_{g_{l-1}} \varepsilon_0^{2q-g_{l-1}} & \cdots & \vdots \\ \vdots & \vdots & \ddots & \vdots \\ \bar{\alpha}_{g_0} \varepsilon_0^{2q-g_0} & \cdots & \vdots \end{vmatrix}.$$

Multiply each of the first l columns of $D'(\alpha, \beta)$ by $e^{i(x-\psi)/2}$.

Multiply each of the last l+1 columns of $D'(\alpha, \beta)$ by $e^{i(\psi-x)/2}$.

The factor $e^{i(\psi + z)/2}$ can be removed from the middle row of the resulting determinant. Changing variables

$$\alpha_i' = e^{i(\psi - \chi)/2}\alpha_i, \qquad \beta_i' = e^{i(\chi - \psi)/2}\beta_i, \quad j \in L,$$

we get $D'(\alpha, \beta) = e^{i\alpha}D(\alpha', \beta')$, where D agrees with D' except that its middle row is

$$\varepsilon_0^q \cdot \cdot \cdot \varepsilon_{q+l}^q$$
.

Recapitulating,

PROPOSITION O. Boundary conditions of selfadjoint type are regular in case (O) if for arbitrary complex numbers α_j , β_j , subject only to the restriction that $|\alpha_j| + |\beta_j| > 0$, $j \in L$, we have $D(\alpha, \beta) \neq 0$.

Proof. The inequality $D(\alpha, \beta) \neq 0$ implies that $D(\alpha', \beta') \neq 0$, and since α_j, β_j are arbitrary, that $D(\beta', \alpha') \neq 0$. Corollary 2 now applies.

In closing this section, we shift back to case (E), where we will remain until the conclusion. The theorem of this paper will be established when Lemma IV is proved. The forthcoming Lemma III will be auxiliary to this proof.

- III. A class of positive definite matrices. In this section $C_k(A)$ will denote the kth compound of the matrix A, with the property that $C_k(AB) = C_k(A)C_k(B)$. [6, pp. 16-17].
- 1. A brief review of positive definiteness. The hermitian matrix $A = (a_{jk})_{j,k=1}^n$ is p.d. if $(Ax, x) = \sum_{j,k=1}^n a_{jk} x_j \bar{x}_k > 0$ for all nonzero vectors $x = (x_1, \dots, x_n)$. It is a consequence that A is p.d. if and only if all its principal minors (the ones chosen from the same rows as columns) are positive, or if and only if $A = B^2$ for some nonsingular hermitian B. If A and B are hermitian and $a_{jk} = d_j d'_k b_{jk}$, where d_j , d'_k are positive, j, $k = 1, \dots, n$, then A is p.d. if and only if B is.

If A is p.d. then so are $A \otimes A$ and $C_k(A)$, $k = 1, \ldots, n$. In fact $B^2 \otimes B^2 = (B \otimes B)^2$ and $C_k(B^2) = (C_k(B))^2$. The matrix of squares $(a_{jk}^2)_{j,k=1}^n$, being a principal submatrix of $A \otimes A$, is p.d. also.

2. A class of positive definite matrices. We require a notation for matrices indexed by sets somewhat more general than the natural numbers. Let Q be the ordered set $\{0, 1, \ldots, q-1\}$. For $r=1, \ldots, q$, Q_r is the set of all strictly increasing sequences in Q of length r; ρ and σ will always be elements of Q_r . Order Q_r by setting $\rho > \sigma$ if the first integer in ρ which differs from the corresponding integer in σ is greater than that integer. Let $(a_{\rho\sigma})_{\rho,\sigma\in Q_r}$ be the $\binom{q}{r} \times \binom{q}{r}$ matrix whose elements $a_{\rho\sigma}$ are indexed by Q_r with the given ordering. Thus $(a_{jk})_{j,k=1}^{q-1}$, indexed by Q_1 , is the usual notation for a $q \times q$ matrix. Its $r \times r$ submatrix with row and column indices in ρ and σ respectively, is written $(a_{jk})_{j\in\rho,k\in\sigma}$.

LEMMA III. The matrices P_r are positive definite for r=1, 2, ..., q, where

$$P_r = \left(\prod_{j \in \rho; k \in \sigma} \operatorname{cosec} (x_j + x_k) \right)_{\rho, \sigma \in Q_r}$$

and $0 < x_0 < x_1 < \cdots < x_{g-1} < \pi/2$.

Proof. First observe that each P_r is a symmetric matrix with positive elements. Now fix $r, 2 \le r \le q$. We define the positive quantities

$$t_j = \tan x_j, \quad j \in Q, \qquad T_\rho = \prod_{j' > j; j, j' \in \rho} (t_{j'} - t_j),$$

$$C_\rho = \prod_{j \in \rho} \cos x_j, \qquad y_{\rho\sigma} = \prod_{j \in \rho; k \in \sigma} \operatorname{cosec}(x_j + x_k).$$

Let $p_{\rho\sigma}$ be the minor of P_1 with row and column indices in ρ and σ . Then

$$p_{\rho\sigma} = \det\left(\csc\left(x_j + x_k\right)\right)_{j \in \rho, k \in \sigma}$$

$$= \det\left(\sec x_j \sec x_k(t_j + t_k)^{-1}\right)_{j \in \rho, k \in \sigma}$$

$$= C_\rho^{-1} C_\sigma^{-1} \det\left((t_j + t_k)^{-1}\right)_{j \in \rho, k \in \sigma}$$

$$= C_\rho^{-1} C_\sigma^{-1} T_\rho T_\sigma \prod_{j \in \rho; k \in \sigma} (t_j + t_k)^{-1}$$

$$= C_\rho^{\tau - 1} C_\sigma^{\tau - 1} T_\rho T_\sigma y_{\rho\sigma}.$$

The second equation is a consequence of the identity $\csc(x+y) = \sec x \sec y$. $(\tan x + \tan y)^{-1}$, and the fourth holds by a lemma of Cauchy [10, Volume 2, pp. 98, 299]. It follows that $p_{\rho\sigma}$ is positive and that

$$(1) P_r = D_r C_r(P_1) D_r.$$

Here D_r is a diagonal matrix with positive diagonal elements $C_\rho^{1-r}T_\rho^{-1}$, $\rho \in Q_r$. In conclusion, P_1 is p.d. (it is even totally definite). Thus $C_r(P_1)$ is p.d., and from (1), P_r is p.d., $r=2,\ldots,q$. We remark that this proof generalizes easily to show that if

- (i) g and h are positive functions defined for the sequence $x_0 < x_1 < \cdots < x_{q-1}$ and h is strictly increasing.
 - (ii) f satisfies the functional equation $f(x+y)=g(x)g(y)[h(x)+h(y)]^{-1}$ then

$$\left(\prod_{j\in\rho;\,k\in\sigma}f(x_j+x_k)\right)_{\rho,\,\sigma\in\mathcal{Q}_r}$$

is p.d. for $r=1,\ldots,q$.

IV. Proof that $D \neq 0$.

1. Some facts about the determinant δ . Let δ denote the $2l \times 2l$ determinant obtained from D by setting $\alpha_i = \beta_i = 1, j \in L$.

The first l columns of δ will be called the left columns and the last l the right columns. Likewise, the first and last l rows will be top and bottom. Let $g_0 \le j_0 < j_1 < \cdots < j_{l-1} \le 2q-1-g_0$ and let

$$M = \begin{vmatrix} \varepsilon_q^{j_0} & \cdots & \varepsilon_{q+l-1}^{j_0} \\ \vdots & & \vdots \\ \varepsilon_{n}^{j_{l-1}} & \cdots & \varepsilon_{n+l-1}^{j_{l-1}} \end{vmatrix}$$

be a minor of δ from the right. Then

(1)
$$M = (-1)^{j_0 + \cdots + j_{l-1}} \begin{vmatrix} \varepsilon_0^{j_0} & \cdots & \varepsilon_{l-1}^{j_0} \\ \vdots & & \vdots \\ \varepsilon_0^{j_{l-1}} & \cdots & \varepsilon_{l-1}^{j_{l-1}} \end{vmatrix}.$$

The latter determinant, to be called m, is the left minor of δ chosen from the same rows as M. By the relation $e_s^s = e_s^t$, m is the Vandermonde

(2)
$$m = \prod_{0 \le k < k' \le l-1} (\varepsilon_{j_{k'}} - \varepsilon_{j_k}).$$

In evaluating this product, we will make use of the identity

(3)
$$\varepsilon_t - \varepsilon_s = 2i \sin \frac{\pi}{2q} (t-s) e^{(i\pi/2q)(t+s)}.$$

Let d_1 , d_2 , d_3 , d_4 be the 'corner' minors of δ :

$$\delta = \begin{vmatrix} d_1 & d_2 \\ -d_3 & d_4 \end{vmatrix}$$

so that

$$d_1 = \begin{vmatrix} \varepsilon_0^{g_0} & \cdots & \varepsilon_{l-1}^{g_0} \\ \vdots & & \vdots \\ \varepsilon_0^{g_{l-1}} & \cdots & \varepsilon_{l-1}^{g_{l-1}} \end{vmatrix}$$
, etc.

Then writing $G = \sum_{j \in L} j$ and $S = \prod_{j' > j; j, j' \in L} \sin(\pi/2q)(j'-j)$, we get from (1), (2) and (3) that

$$d_1 = \prod_{j'>j; j,j' \in L} (\varepsilon_{j'} - \varepsilon_j) = (2i)^{l(l-1)/2} Se^{(i\pi/2q)(l-1)G}$$

and

$$d_4 = (-1)^{l+G} (2i)^{l(l-1)/2} Se^{(i\pi/2q)(1-l)(l+G)}.$$

Thus

$$d_1d_4=(-1)^{l+G}(2i)^{l(l-1)}S^2e^{(i\pi/2q)l(1-l)},$$

and employing (1),

$$d_2d_3 = (-1)^l d_1d_4.$$

2. Proof that $D \neq 0$. For r = 1, 2, ..., l we define L_r to be the set of all strictly increasing sequences in L of length r, and give L_r the ordering of Q_r (III.2). Then $L_r \subset Q_r$ and $(y_{\rho\sigma})_{\rho,\sigma\in L_r}$, a principal submatrix of P_r , is p.d. This will be basic in the proof to follow.

LEMMA IV. For arbitrary complex numbers α_j , β_j , subject only to the restriction that $|\alpha_j| + |\beta_j| > 0$, $j \in L$, the $2l \times 2l$ determinant D of II.3 does not vanish.

Proof. The mapping $\hat{}: L \to \{0, 1, ..., l-1\}$ will send $j \in L$ onto \hat{j} determined uniquely by $j = g_{\hat{j}}$. For $j \in L$ we define $n_j = j + \hat{j}$, $N = \sum_{j \in L} n_j$ and $\hat{\beta}_j = (-1)^{n_j} \beta_j$.

The integer r will now be fixed, $1 \le r \le l-1$. The symbols ρ and σ will always denote elements of L_r . We make the following definitions:

$$N_{\rho} = \sum_{j \in \rho} n_{j},$$

$$x_{j} = (j + \frac{1}{2})\pi/2q, \quad j \in \rho,$$

$$S_{\rho} = \prod_{j' > j} \sin(x_{j'} - x_{j}),$$

$$S'_{\rho} = S_{\rho} \prod_{j \in \rho; k \in L} \sin(x_{j} + x_{k}),$$

$$c = (-1)^{l+N} 2^{l(l-1)} e^{(t\pi/2q)l(1-l)}.$$

A second mapping $\tilde{c}: L_r \to L_{l-r}$ sends ρ onto $\tilde{\rho} \in L_{l-r}$ consisting of those l-r integers in L disjoint from ρ . The minor $m_{\rho\sigma}$ is defined to be

$$\begin{vmatrix} \varepsilon_0^{j_0} & \cdots & \varepsilon_{l-1}^{j_0} \\ \vdots & & \vdots \\ \varepsilon_0^{j_{r-1}} & \cdots & \varepsilon_{l-1}^{j_{r-1}} \\ \varepsilon_0^{2q-1-k_{l-r-1}} & \cdots & \varepsilon_{l-1}^{2q-1-k_{l-r-1}} \\ \vdots & & \vdots \\ \varepsilon_0^{2q-1-k_0} & \cdots & \varepsilon_{l-1}^{2q-1-k_0} \end{vmatrix}$$

where $\rho = \{j_0, \ldots, j_{\tau-1}\}$ and $\tilde{\sigma} = \{k_0, \ldots, k_{l-\tau-1}\}$. In words, $m_{\rho\sigma}$ is the left minor of δ chosen from the rows indicated by ρ in the top half and the ones indicated by $\tilde{\sigma}$, in reverse order, in the bottom half. Let $M_{\rho\sigma}$ be the cofactor of $m_{\rho\sigma}$ in δ . We then have the formula

$$(4) m_{\rho\sigma}M_{\rho\sigma} = m_{\sigma\rho}M_{\sigma\rho} = (-1)^{N_{\rho}+N_{\sigma}}S'_{\rho}S_{\tilde{\rho}}S'_{\sigma}S_{\tilde{\sigma}}y_{\rho\sigma}^{2}.$$

To prove (4) define

$$y_{\rho\tilde{\sigma}} = \prod_{j\in\rho; k\in\tilde{\sigma}} \operatorname{cosec}(x_j + x_k) = y_{\tilde{\sigma}\rho}.$$

Since $y_{\rho\tilde{\sigma}}y_{\rho\sigma} = \prod_{i \in \rho: k \in L} \operatorname{cosec}(x_i + x_k)$, we get

$$S_{\rho}y_{\rho\bar{\sigma}}^{-1} = S_{\rho}'y_{\rho\sigma}.$$

Now using (2) and (3),

$$m_{og} = (2i)^{l(l-1)/2} \Pi_1 \Pi_2 \Pi_3$$

where

$$\Pi_{1} = \prod_{j'>j; j,j'\in\rho} \sin(x_{j'}-x_{j}) e^{(i\pi/2q)(j+j')} = S_{\rho} \exp\left[\frac{i\pi}{2q}(r-1)\sum_{j\in\rho}j\right],$$

$$\Pi_{2} = \prod_{k'>k; k,k'\in\tilde{\sigma}} \sin(x_{k'}-x_{k}) e^{(i\pi/2q)(4q-2-k-k')}$$

$$= S_{\tilde{\sigma}} \exp\left[\frac{i\pi}{2q}(r+1-l)\sum_{k\in\tilde{\sigma}}(k+1)\right],$$

$$\Pi_{3} = \prod_{j\in\rho; k\in\tilde{\sigma}} \sin(x_{j}+x_{k}) e^{(i\pi/2q)(2q-1-k+j)}$$

$$= (-1)^{r(l-r)}y_{\rho\tilde{\sigma}}^{-1} \exp\frac{i\pi}{2q}\left[(l-r)\sum_{j\in\rho}j-r\sum_{k\in\tilde{\sigma}}(k+1)\right].$$

Empty products, which arise if r=1 or l-1, will be interpreted as equal to one. Multiplying the three partial products,

$$m_{\rho\sigma} = (-1)^{r(l-r)} (2i)^{l(l-1)/2} S_{\rho} S_{\tilde{\sigma}} y_{\rho\tilde{\sigma}}^{-1} \exp(i\pi/2q)(l-1) \left[\sum_{i \in \tilde{\sigma}} j - \sum_{k=1}^{r} (k+1) \right].$$

The cofactor $M_{\rho\sigma}$ is by definition the signed minor of δ from the right, the rows indicated by $\bar{\rho}$ in the top and by σ , in reverse order, in the bottom i.e. the rows omitted by $m_{\rho\sigma}$. The exponent of -1 in the sign, call it s_1 , carried by $M_{\rho\sigma}$ is the sum of the row and column indices in $m_{\rho\sigma}$ (counting from 1 to 2*l*) and this sum is

$$\sum_{j \in o} (j+1) + \sum_{k \in \bar{o}} (2l - \hat{k}) + l(l+1)/2.$$

By (1) the unsigned minor $s_1 M_{\rho\sigma}$ is equal to the factor $s_2 = (-1)^{\sum_{j \in \tilde{\rho}} j + \sum_{k \in \sigma} (k+1)}$ times the left minor $m_{\tilde{\rho}\tilde{\sigma}}$ chosen from the same rows.

The exponent of -1 in s_1s_2 is

$$\sum_{j\in\rho} (j+1) + \sum_{j\in\bar{\rho}} j + \sum_{k\in\sigma} (k+1) + \sum_{k\in\bar{\sigma}} (2l-\hat{k}) + l(l+1)/2$$

$$\equiv \sum_{j\in\rho} n_j + \sum_{k\in\sigma} n_k + \sum_{k\in L} n_k + l(l+1)/2 \pmod{2}.$$

Thus

$$s_1 s_2 = (-1)^{N_\rho + N_\sigma + N + l(l+1)/2}$$

and

$$m_{\rho\sigma}M_{\rho\sigma} = s_1 s_2 m_{\rho\sigma} m_{\tilde{\rho}\tilde{\sigma}}$$

$$= c(-1)^{N_{\rho} + N_{\sigma}} S_{\rho} S_{\tilde{\sigma}} S_{\sigma} S_{\tilde{\sigma}} y_{\rho\tilde{\sigma}}^{-1} y_{\sigma\tilde{\rho}}^{-1}$$

$$= c(-1)^{N_{\rho} + N_{\sigma}} S_{\rho}' S_{\tilde{\rho}} S_{\sigma}' S_{\tilde{\sigma}} y_{\rho\sigma}^{2} \quad \text{(by (5))}$$

$$= m_{\sigma}M_{\sigma\sigma}$$

Equation (4) is verified.

If $\gamma = (\gamma_0, \gamma_1, \dots, \gamma_{l-1})$ is any vector of l complex components let $\gamma^{\rho} = \prod_{j \in \rho} \gamma_j$, and $\bar{\gamma}^{\rho} = \prod_{j \in \rho} \bar{\gamma}_j$. Also, let $\hat{\beta}^{\rho} = \prod_{j \in \rho} \hat{\beta}_j$. In this notation

$$\begin{vmatrix} \beta_{j_0} \varepsilon_0^{j_0} & \cdots & \beta_{j_0} \varepsilon_{l-1}^{j_0} \\ \vdots & & \vdots \\ \beta_{j_{r-1}} \varepsilon_0^{j_{r-1}} & \cdots & \beta_{j_{r-1}} \varepsilon_{l-1}^{j_{r-1}} \\ \bar{\alpha}_{k_{l-r-1}} \varepsilon_0^{2q-1-k_{l-r-1}} & \cdots & \bar{\alpha}_{k_{l-r-1}} \varepsilon_{l-1}^{2q-1-k_{l-r-1}} \\ \vdots & & \vdots \\ \bar{\alpha}_{k_0} \varepsilon_0^{2q-1-k_0} & \cdots & \bar{\alpha}_{k_0} \varepsilon_{l-1}^{2q-1-k_0} \end{vmatrix} = \beta^{\rho} \bar{\alpha}^{\bar{\beta}} m_{\rho\sigma}.$$

The cofactor of this term in the Laplace expansion of D from the first l columns is $\alpha^{\bar{\rho}} \bar{\beta}^{\sigma} M_{\rho \sigma}$ (Figure 3).

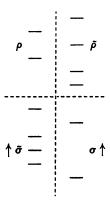


FIGURE 3

Now $\beta^{\rho} = \prod_{i \in \rho} \beta_i = \prod_{j \in \rho} (-1)^{n_j} \hat{\beta}_j = (-1)^{N_{\rho}} \hat{\beta}^{\rho}$. Thus

$$\beta^{\rho}\alpha^{\tilde{\rho}}\bar{\beta}^{\sigma}\bar{\alpha}^{\tilde{\sigma}}m_{\rho\sigma}M_{\rho\sigma} = \hat{\beta}^{\rho}\alpha^{\tilde{\rho}}\hat{\beta}^{\sigma}\bar{\alpha}^{\tilde{\sigma}}(cS_{\rho}'S_{\tilde{\sigma}}S_{\sigma}'S_{\tilde{\sigma}}v_{\rho\sigma}^{2}).$$

The sum

$$\sum_{\rho,\sigma\in L_\tau}\beta^\rho\alpha^{\tilde{\rho}}\bar{\beta}^\sigma\bar{\alpha}^{\tilde{\sigma}}m_{\rho\sigma}M_{\rho\sigma}$$

contains just those terms in the expansion of D obtained by taking r of the rows in the left minor from the top. It is a quadratic form $(A_r\gamma_r, \gamma_r)$ in the variables $\gamma_r = (\hat{\beta}^\rho \alpha^{\bar{\rho}})_{\rho \in L_r}$ with coefficient matrix $A_r = c(S'_\rho S_{\bar{\rho}} S'_\sigma S_{\bar{\rho}} y^2_{\rho\sigma})_{\rho,\sigma \in L_r}$. As remarked in the first paragraph of this section, $(y_{\rho\sigma})_{\rho,\sigma \in L_r}$ is p.d. Since S'_ρ and $S_{\bar{\rho}}$ are positive, $c^{-1}A_r$ is also p.d. This holds for $r = 1, 2, \ldots, l-1$.

Let $\gamma_0 = \prod_{j \in L} \alpha_j \bar{\alpha}_j$, $\gamma_l = \prod_{j \in L} \beta_j \bar{\beta}_j = \prod_{j \in L} \hat{\beta}_j \bar{\beta}_j$. The matrices $A_0 = A_l$ accompanying γ_0 , γ_l in the Laplace expansion each have one element $d_1 d_4 = (-1)^l d_2 d_3$. Since

$$N = \sum_{i \in L} n_i = \sum_{i \in L} j + \sum_{k=0}^{l-1} k = G + l(l-1)/2$$

we conclude that $d_1d_4 = cS^2$.

We have shown that

$$D = \sum_{r=0}^{l} (A_r \gamma_r, \gamma_r)$$

and that $c^{-1}A_{\tau}$ is p.d., r=0, 1, ..., l. By hypothesis $|\alpha_j| + |\beta_j| > 0$, $j \in L$. Hence there exists an $r, 0 \le r \le l$, and a $\rho \in L_{\tau}$ such that $\hat{\beta}^{\rho} \alpha^{\tilde{\rho}} \ne 0$. The vector γ_{τ} is not zero, and the lemma is established.

It is possible to furnish a shorter proof if l=q [13].

3. An illustration. In Table 1 we illustrate this lemma for the case q=3.

Table 1

L	D D
{1, 2}	$-e^{2i\pi/3}[\beta_1\beta_2 ^2+ \alpha_1\alpha_2 ^2+ \beta_1\alpha_2 ^2+ \beta_2\alpha_1 ^2+2 \beta_1\alpha_2+\beta_2\alpha_1 ^2]$
{0, 2}	$e^{2i\pi/3}[3(\beta_0\beta_2 ^2+ \alpha_0\alpha_2 ^2+ \beta_0\alpha_2 ^2+ \beta_2\alpha_0 ^2)+ \beta_0\alpha_2-\beta_2\alpha_0 ^2]$
{0, 1}	$-e^{2i\pi/3}[\beta_0\beta_1 ^2+ \alpha_0\alpha_1 ^2+ \beta_0\alpha_1 ^2+ \beta_1\alpha_0 ^2+2 \beta_0\alpha_1+\beta_1\alpha_0 ^2]$
{0, 1, 2}	

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University of California, Berkeley, California