

# AN EXTENSION OF A THEOREM OF G. SZEGÖ AND ITS APPLICATION TO THE STUDY OF STOCHASTIC PROCESSES<sup>(1)</sup>

BY

ULF GRENANDER AND MURRAY ROSENBLATT

1. **Introduction.** In this paper we study minimum problems associated with quadratic forms

$$(1.1) \quad Q_n = c' M^{(n)} c$$

where  $c$  is a column vector with components  $c_0, c_1, \dots, c_n$  and  $M^{(n)}$  is a Hermitian matrix with the elements

$$m_{p,q}^{(n)} = \int_{-\pi}^{\pi} e^{i(p-q)\lambda} f(\lambda) d\lambda, \quad p, q = 0, 1, \dots, n.$$

We denote the conjugate of the transpose of a matrix  $A$  by  $A'$ . Here  $f(\lambda)$  is a nonnegative integrable function in  $(-\pi, \pi]$ . We shall define  $f(\lambda)$  with period  $2\pi$  on the real axis. Some of these minimum problems arise in the theory of stationary stochastic processes. These applications will be discussed in §5 [3].

Szegö [6] has studied the minimum  $\mu_n$  of  $Q_n$  subject to the restraint

$$P_n(\alpha) = 1$$

where

$$P_n(w) = \sum_{r=0}^n c_r w^r.$$

He has shown that if  $|\alpha| < 1$ , the limit  $\mu$  of  $\mu_n$  as  $n \rightarrow \infty$  is positive if and only if

$$\int_{-\pi}^{\pi} \log f(\lambda) d\lambda > -\infty.$$

Then we can write the formal Fourier expansion

$$\log f(\lambda) \sim k_0 + 2 \sum_{\nu=1}^{\infty} (k_{\nu} \cos \nu\lambda + l_{\nu} \sin \nu\lambda).$$

Putting

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$$g(w) = \frac{k_0}{2} + \sum_{r=1}^{\infty} (k_r - il_r)w^r$$

and

$$D(w) = e^{g(w)},$$

Szegö has shown that

$$\mu = 2\pi |D(\alpha)|^2(1 - |\alpha|^2).$$

We are going to study the minimization of the quadratic form  $Q_n$  with restraints

$$(1.2) \quad C: P^{(k)}(\alpha_j) = \beta_j^k, \quad k = 0, 1, \dots, n_j, j = 1, \dots, m,$$

where the  $\alpha_j$  are different points in the closed unit circle and where the  $\beta_j^k$  do not all vanish. The results of the paper are valid with appropriate modification when the restraints are of the form

$$C^*: P^{(k)}(\alpha_j) = \beta_j^k, \quad k \in S_j, j = 1, \dots, m,$$

where  $S_j$  is a finite set of nonnegative integers. We have restricted ourselves to restraints of the form (1.2) in order to avoid excessive notation.

**2. Conditions inside the unit circle.** We order the pairs  $(j, k)$  according to increasing  $j$  and for fixed  $j$  according to increasing  $k$ ; let  $r$  be the numbering index of these pairs,  $r = 1, 2, \dots, N = \sum_{j=1}^m (n_j + 1)$ . Defining the inner product of two polynomials  $g(w), h(w)$  as

$$(g, h) = \int_{-\pi}^{\pi} g(e^{i\lambda}) [h(e^{i\lambda})]^* f(\lambda) d\lambda^{(2)},$$

we introduce the orthonormal polynomials  $\phi_\nu(w)$ ,  $\nu = 0, 1, \dots$ , obtained by the Gram-Schmidt procedure from  $1, w, w^2, \dots$  [6]. Then we can write

$$P_n(w) = \sum_{\nu=0}^n d_\nu \phi_\nu(w)$$

so that

$$Q_n = \sum_{\nu=0}^n |d_\nu|^2.$$

**THEOREM 1<sup>(3)</sup>.**

$$(2.1) \quad \mu_n = \beta'(H_n')^{-1}\beta, \quad n > N,$$

(<sup>2</sup>)  $[\dots]^*$  denotes the conjugate of  $[\dots]$ .

(<sup>3</sup>) We thank the referee for suggesting the simple proof of Theorem 1 given above.

where  $\beta$  is a column vector with the  $N$  components  $\beta_r = \beta_j^k$  and  $H_n$  is a nonsingular  $N \times N$  matrix with elements

$$h_{r,s} = \sum_{\nu=0}^n \phi_{\nu}^{(k)}(\alpha_j) \phi_{\nu}^{(k')}(\alpha_{j'}),$$

$r \leftrightarrow (j, k), s \leftrightarrow (j', k')$ .

If  $\log f(\lambda)$  is integrable and all the restraints (1.2) are at points  $\alpha_j$  inside the unit circle (we call such a set of restraints  $C_i$ ), then

$$(2.2) \quad \mu = \lim_{n \rightarrow \infty} \mu_n = \beta'(H')^{-1}\beta$$

where  $H$  is a nonsingular  $N \times N$  matrix with elements

$$h_{r,s} = \frac{1}{2\pi} \left( \frac{\partial^k}{\partial x^k} \frac{\partial^{k'}}{\partial y^{*k'}} \frac{1}{1 - xy^*} \frac{1}{D(x)[D(y)]^*} \right), \quad x = \alpha_j, y = \alpha_{j'}.$$

**Proof.** We have to minimize  $Q_n = \sum_{\nu=0}^n |d_{\nu}|^2, n \geq N$ , with the restraint

$$\sum_{\nu=0}^n d_{\nu} \phi_{\nu}^{(k)}(\alpha_j) = \beta_j^k, \quad \text{or} \quad \sum_{\nu=0}^n d_{\nu} l_{\nu r} = \beta_r, \quad r = 1, 2, \dots, N.$$

We introduce the column vectors  $d = \{d_{\nu}^*\}, l_r = \{l_{\nu r}\} = \{\phi_{\nu}^{(k)}(\alpha_j)\}$  in  $n$ -space so that the restraints have the form  $d'l_r = \beta_r$  and  $|d|^2$  has to be minimized. The vectors  $l_r$  are linearly independent since

$$\sum l_j \phi_{\nu}^{(k)}(\alpha_j) = 0, \quad \nu = 0, 1, \dots, n,$$

would mean that

$$\sum_{j,k} l_j f^{(k)}(\alpha_j) = 0$$

for any polynomial  $f(z)$  of degree  $n$ . By proper choice of  $f$  we find  $l_j^* = 0$ .

Now the vectors  $l'_r$  span a linear manifold

$$\sum_{r=1}^N \lambda_r l'_r$$

and the projection of  $d$  onto this furnishes the minimum of  $|d|$ . So assume  $d$  has the latter form. The restraints are now

$$\sum_{r=1}^N \lambda_r (l_s, l_r) = \beta_s, \quad s = 1, 2, \dots, N.$$

The Hermitian matrix  $H_n = [(l_s, l_r)]$  is positive definite. Introducing the column vector  $\lambda = \{\lambda_r\}$ , we can write these equations as follows:  $H_n \lambda = \beta$  so that  $\lambda = H_n^{-1} \beta$ . The minimum in question is

$$\left| \sum_{r=1}^N \lambda_r l_r' \right|^2 = \sum_{r,s} \lambda_r \bar{\lambda}_s (l_s, l_r) = \lambda' H_n \lambda = \beta' (H_n')^{-1} H_n H_n^{-1} \beta = \beta' (H_n')^{-1} \beta.$$

Substituting we obtain for the elements of  $H_n$ :

$$(l_r, l_s) = \sum_{\nu=0}^n \phi_\nu^{(k)}(\alpha_j) [\phi^{(k')}(\alpha_{j'})]^*.$$

But

$$\lim_{n \rightarrow \infty} \sum_{\nu=0}^n \phi_\nu(x) [\phi_\nu(y)]^* = \frac{1}{2\pi} \frac{1}{1 - xy^*} \frac{1}{D(x)[D(y)]^*}$$

uniformly for  $|x|, |y| \leq r < 1$  [6, Satz XXXI]. From this it easily follows that

$$\begin{aligned} \lim_{n \rightarrow \infty} \sum_{\nu=0}^n \phi_\nu^{(k)}(\alpha_j) [\phi_\nu^{(k')}(\alpha_{j'})]^* \\ = \frac{1}{2\pi} \left( \frac{\partial^k}{\partial x^k} \frac{\partial^{k'}}{\partial y^{*k'}} \frac{1}{1 - xy^*} \frac{1}{D(x)[D(y)]^*} \right), \quad x = \alpha_j, y = \alpha_{j'}, \end{aligned}$$

and we have (2.2).

3. **Conditions on the unit circle.** We now consider the restraints (1.2) at points  $\alpha_j$  on the unit circle and call such a set of restraints  $C_b$ . This case differs considerably from that just treated. Here we get  $\mu = 0$  and are mainly interested in the principal term of  $\mu_n$  as  $n \rightarrow \infty$ . To study this we have to introduce certain regularity conditions on  $f(\lambda)$ .

THEOREM 2. *Let*

$$(3.1) \quad f(\lambda) = g(\lambda) \prod_{\nu=1}^{\infty} |e^{i\lambda} - e^{i\theta_\nu}|^{2l_\nu}, \quad -\pi < \theta_\nu \leq \pi,$$

where  $g(\lambda)$  is positive and continuous and the  $l_\nu$  are positive integers. Then

$$(3.2) \quad \begin{aligned} \mu_n &= \frac{2\pi}{n^{2\rho+1}} \sum_{r_j+t_j=\rho} \binom{r_j+t_j}{r_j}^2 \cdot |\beta_j^{r_j}|^2 \cdot d(\rho, \nu_j) \cdot g(\lambda_j) \cdot \prod_{e^{i\theta_\nu} \neq \alpha_j} |e^{i\theta_\nu} - \alpha_j|^{2l_\nu} \\ &+ O\left(\frac{1}{n^{2\rho+2}}\right), \quad \alpha_j = e^{i\lambda_j}, \end{aligned}$$

where  $\rho, r_j, t_j, \nu_j$  and  $d(\rho, \nu_j)$  are defined below.

**Proof.** Choose two trigonometric polynomials  $a(w), b(w)$  in  $w = e^{i\lambda}$  of order  $p$  so that

$$\frac{1}{2\pi} |a(e^{i\lambda})|^2 \leq g(\lambda) \leq \frac{1}{2\pi} |b(e^{i\lambda})|^2, \quad |b(e^{i\lambda})| - |a(e^{i\lambda})| < \epsilon.$$

Let us now consider the case when  $g(\lambda)$  is exactly equal to

$$\frac{1}{2\pi} |a(e^{i\lambda})|^2 = \frac{1}{2\pi} \left| \sum_{r=0}^p a_r e^{ir\lambda} \right|^2.$$

Then we should minimize under conditions  $C_b$

$$\frac{1}{2\pi} \int_{-\pi}^{\pi} |P_n(e^{i\lambda}) p(e^{i\lambda})|^2 d\lambda$$

where

$$p(w) = a(w) \prod_{\nu} (w - z_{\nu})^{l_{\nu}}, \quad z_{\nu} = e^{i\theta_{\nu}}.$$

Let  $q(w) = P_n(w)p(w)$ . The problem can then be rephrased in the following manner. We minimize the integral

$$\frac{1}{2\pi} \int_{-\pi}^{\pi} |q(e^{i\lambda})|^2 d\lambda$$

under the conditions, say  $C_b^*$ , that the conditions  $C_b$  induce on  $q(w)$ . It is clear that the  $C_b^*$  are of a similar form

$$q^{(k)}(w_j) = \eta_j^k, \quad k = 0, 1, \dots, \nu_j, \quad j = 1, 2, \dots, M,$$

where  $w_j$  is one of the  $\alpha_j$ 's or  $z_{\nu}$ 's. The range of  $(j, k)$  is not necessarily the same as in  $C_b$  but we carry out in a similar mapping of  $(j, k)$  onto a single index  $r$ . Then we get from (2.1)

$$(3.3) \quad \mu_n = \eta'(H'_n)^{-1} \eta.$$

To compute  $H_n$  we observe that in the present case  $f(\lambda) = 1/2\pi$  so that  $\phi_r(w) = w^r$ . Hence if  $|x| = |y| = 1$ ,

$$\begin{aligned} & \sum_{r=0}^n \phi_r^{(k)}(x) [\phi_r^{(k')}(y)]^* \\ &= x^{-k} y^{k'} \sum_{r=0}^n \nu(\nu-1) \cdots (\nu-k+1) \nu(\nu-1) \cdots (\nu-k'+1) x^r y^{-r}. \end{aligned}$$

If  $x=y$  we get

$$x^{-k+k'} \frac{n^{k+k'+1}}{k+k'+1} + O(n^{k+k'}).$$

If  $x \neq y$  we use the Abel summation formula and find that the expression is  $O(n^{k+k'})$ . Hence



and

$$q^{(k)}(\alpha_j) = 0, \quad k < r_j.$$

Now let  $\alpha_j = z_\nu$ . Let  $r_j$  be defined as above. Then

$$q^{(k)}(z_\nu) = 0, \quad k < r_j + l_\nu,$$

$$q^{(r_j+l_\nu)}(z_\nu) = \binom{r_j + l_\nu}{r_j} \beta_j^{r_j} a(z_\nu) \prod_{\mu \neq \nu} (z_\nu - z_\mu)^{l_\mu}.$$

Let  $t_j$  be the order of the zero of  $q(w)$  at  $\alpha_j$ . If there is no zero we set  $t_j = 0$  and if  $\alpha_j = z_\nu$ , then  $t_j = l_\nu$ . Let the least  $k$  such that there is an  $\eta_j^k \neq 0$  be called  $\rho$ . It follows from (3.3), (3.4), and (3.5) that the leading term of  $\mu_n$  will consist only of contributions from the conditions with  $k = \rho$ . It is clear from the discussion above that we then get

$$\mu_n = \frac{1}{n^{2\rho+1}} \sum_{r_j+t_j=\rho} \binom{r_j + t_j}{r_j}^2 \cdot |\beta_j^{r_j}|^2 \cdot \left| a(\alpha_j) \prod_{z_\mu \neq \alpha_j} (\alpha_j - z_\mu)^{l_\mu} \right|^2 \cdot d(\rho, \nu_j) + O\left(\frac{1}{n^{2\rho+2}}\right)$$

where

$$d(\rho, n) = \left\{ \frac{1}{\nu + \mu + 1}; \nu, \mu = 0, \dots, n \right\}_{\rho, \rho}^{-1}$$

$$= (2\rho + 1) \prod_{j=0, j \neq \rho}^n \left( \frac{j + \rho + 1}{j - \rho} \right)^2$$

(see [1, p. 177]).

The true value of  $\mu_n$  if  $f(\lambda)$  is given by (3.1) will be between the two values computed for  $a(w)$ ,  $b(w)$ . Letting  $\epsilon$  tend to zero we get the desired result.

REMARK 1. If none of the zeros of  $f(\lambda)$  coincide with any  $\alpha_j$ , then  $\rho$  is simply the smallest  $k$  such that there is a nonzero  $\beta_j^k$ . We then have

$$\mu_n = \frac{2\pi}{n^{2\rho+1}} \sum_j |\beta_j^\rho|^2 \cdot f\left(\frac{\log \alpha_j}{i}\right) \cdot d(\rho, n_j) + O\left(\frac{1}{n^{2\rho+2}}\right).$$

REMARK 2. The case when there are conditions both inside and on the unit circle is now easily handled. If some  $\beta$  in  $C_i$  does not vanish, then  $\mu > 0$  and its value can be computed from (2.2) as if no conditions on the circle had been present, as is easily verified. On the other hand, if all  $\beta$ 's in  $C_i$  vanish,  $\mu_n \rightarrow 0$  and we get its principal term from (3.2) as if  $C_b$  were the only conditions.

Theorem 2 suggests that the error  $\mu'_n$  computed with the minimizing poly-

nomial corresponding to a uniform weight function is of the same order as the error  $\mu_n$  computed with the minimizing polynomial corresponding to weight function  $f(\lambda)$ . The following theorem indicates that this conjecture is essentially true [4].

**THEOREM 3.** *Let  $f(\lambda)$  be a nonnegative continuous function having no zeros in common with the points  $\alpha_j$  of the conditions  $C_b$ . Let*

$$\mu'_n = \int_{-\pi}^{\pi} |P_n(e^{i\lambda})|^2 f(\lambda) d\lambda$$

where  $P_n(w)$  minimizes  $\int_{-\pi}^{\pi} |P_n(e^{i\lambda})|^2 d\lambda$  under conditions  $C_b$ . Then

$$\lim_{n \rightarrow \infty} \frac{\mu'_n}{\mu_n} = 1.$$

**Proof.** Let  $P_n(w) = \sum_0^n \gamma_r w^r$  be the minimizing polynomial under conditions  $C_b$  and the assumption that the spectral density is uniform ( $f(\lambda) = 1$ ). Let

$$\frac{d^k}{dw^k} w^r = \psi_{r,k}(w).$$

Then

$$\frac{1}{2\pi} \int_{-\pi}^{\pi} |P_n(e^{i\lambda})|^2 d\lambda = \sum_0^n |\gamma_r|^2$$

where

$$\gamma_r = \sum_{j,k} \lambda_j^k [\psi_{r,k}(\alpha_j)]^*$$

or

$$\gamma = \psi \lambda$$

where

$$\psi = \left\{ \begin{array}{cccc} [\psi_{0,0}(\alpha_1)]^* & [\psi_{0,1}(\alpha_1)]^* & \cdots & [\psi_{0,n_1}(\alpha_1)]^* & [\psi_{0,0}(\alpha_2)]^* & \cdots \\ [\psi_{1,0}(\alpha_1)]^* & [\psi_{1,1}(\alpha_1)]^* & \cdots & [\psi_{1,n_1}(\alpha_1)]^* & [\psi_{1,0}(\alpha_2)]^* & \cdots \\ \dots & \dots & \dots & \dots & \dots & \dots \\ [\psi_{n,0}(\alpha_1)]^* & [\psi_{n,1}(\alpha_1)]^* & \cdots & [\psi_{n,n_1}(\alpha_1)]^* & [\psi_{n,0}(\alpha_2)]^* & \cdots \end{array} \right\}$$

and  $\lambda = H_n^{-1} \beta$ . As before we factor  $H_n$  so that  $H_n = \mathcal{D}_n \Lambda_n \mathcal{D}_n$ . Then

$$\begin{aligned} \mu'_n &= \sum \gamma_r m_{r,\mu}^{(n)} \gamma_\mu^* = \gamma' M^{(n)} \gamma = \lambda' \psi' M^{(n)} \psi \lambda = \beta' (H_n')^{-1} \psi' M^{(n)} \psi H_n^{-1} \beta \\ &= \beta' \mathcal{D}_n^{-1} (\Lambda_n')^{-1} \mathcal{D}_n^{-1} \psi' M^{(n)} \mathcal{D}_n^{-1} \Lambda_n^{-1} \mathcal{D}_n^{-1} \beta. \end{aligned}$$



$$|\eta_{n,k}(e^{i(\lambda-\lambda_j)})| \leq n^{1/2}.$$

Hence if  $\lambda_j \neq \lambda_{j'}$ , expression (3.6) converges to zero. Now consider  $\alpha_j = \alpha_{j'}$ ,  $= e^{i\lambda_j}$ . Then (3.6) can be rewritten as

$$1/n^{(k+k'+1)} e^{i(k'-k)\lambda_j} \int_{-\pi}^{\pi} \sum_{\nu, \mu=0}^n \nu(\nu-1) \cdots (\nu-k+1) \mu(\mu-1) \cdots (\mu-k'+1) e^{i(\nu-\mu)(\lambda-\lambda_j)} f(\lambda) d\lambda.$$

Let

$$s_k(\lambda) = \frac{1}{n^{k+1/2}} \sum_0^n e^{i(\lambda-\lambda_j)\nu} \nu(\nu-1) \cdots (\nu-k+1).$$

Consider

$$K_n(\lambda) = \sum_{l,k=1}^{n_j} y_l s_l(\lambda) [y_k s_k(\lambda)]^* = \left| \sum y_k s_k(\lambda) \right|^2 \geq 0.$$

One can verify that

1.  $K_n(\lambda) \rightarrow 0$  uniformly in  $\lambda$  if  $|\lambda - \lambda_j| > \epsilon$ ,
2.  $\lim_{n \rightarrow \infty} \int_{-\pi}^{\pi} K_n(\lambda) d\lambda = 2\pi \sum \frac{y_k y_l^*}{k+l+1}$ .

Hence (see [5, p. 49])

$$\lim_{n \rightarrow \infty} \int_{-\pi}^{\pi} s_l(\lambda) s_k(\lambda) f(\lambda) d\lambda = \frac{2\pi f(\lambda_j)}{l+k+1}.$$

But (3.7) then follows immediately. Now

$$\begin{aligned} \mu'_n &= \beta' \mathcal{D}_n^{-1} ((\Lambda')^{-1} + o(1)) (f\Lambda + o(1)) (\Lambda^{-1} + o(1)) \mathcal{D}_n^{-1} \beta \\ &= \beta' \mathcal{D}_n^{-1} ((\Lambda')^{-1} + o(1)) f \mathcal{D}_n^{-1} \beta. \end{aligned}$$

The theorem then follows.

4. **The approach of  $\mu_n - \mu$  to zero under the restraint  $P_n(0) = 1$ .** This problem is of interest in the theory of stochastic processes. Moreover, it does give some insight into the more general problem where the restraints are of the type  $C_i$ . Let  $\delta_n = \mu_n - \mu$ .

**THEOREM 4.** *The decrease of  $\delta_n$  to zero is at least exponential if and only if (1)  $f(\lambda)$  coincides almost everywhere with a function  $g(\lambda)$  that is analytic for all real  $\lambda$  and (2)  $g(\lambda)$  has no zeros.*

**Proof.** Assume that (1) and (2) are satisfied. Then the function

$$\phi(w) = f(\lambda), \quad w = e^{i\lambda}, \quad -\infty < \lambda < \infty,$$

where we have chosen one determination of the logarithmic function, can be analytically extended to an annular region  $\rho_1 < |w| < \rho_2$ ,  $\rho_1 < 1 < \rho_2$ . In this region we can then represent  $\log \phi(w)$  as a convergent Laurent series

$$\log \phi(w) = \sum_{\nu=-\infty}^{\infty} \gamma_{\nu} w^{\nu}.$$

We then note that

$$D(w) = \exp \left\{ \frac{\gamma_0}{2} + \sum_1^{\infty} \gamma_{\nu}^* w^{\nu} \right\}$$

and hence  $D(w)$  is analytic in the closed region  $|w| \leq 1$  and has no zeros in this region. One can define the inner product of two functions  $g(w)$ ,  $h(w)$  such that  $g(w)D(w)$ ,  $h(w)D(w) \in H_2$  as in 2. Then  $\|g\|^2 = (g, g)$  and the set of functions  $g(w)$  such that  $g(w)D(w) \in H_2$  is a Hilbert space. Now

$$\begin{aligned} \mu_n - \mu &\leq \|s_n(w)\|^2 - \|D(0)/D(w)\|^2 \\ &= (\|s_n(w)\| + \|D(0)/D(w)\|)(\|s_n(w)\| - \|D(0)/D(w)\|) \\ &\leq K_1 \left\| s_n(w) - \frac{D(0)}{D(w)} \right\| \leq K_2 \left( \sum_{n+1}^{\infty} |d_{\nu}|^2 \right)^{1/2} \end{aligned}$$

where  $s_n(w) = \sum_0^n d_{\nu} w^{\nu}$  is the  $n$ th partial sum of the Taylor expansion of  $D(0)/D(w)$ . But  $|d_{\nu}| < d^{\nu} < 1$  so that  $\delta_n \leq Kd^{n/2}$ ,  $0 < d < 1$ .

Now assume  $\delta_n \leq Kd^n$ ,  $0 < d < 1$ . Then

$$|\phi_{\nu}(0)| < K_3 d^{n/2}.$$

However,  $|\phi_{\nu}(w)| < K_4 |w|^{\nu}$  on  $|w| = 1 + \epsilon$ ,  $\epsilon > 0$  (see [6, Satz XXXII]). If  $1 + \epsilon < 1/d^{1/2}$  we have uniform convergence of

$$2\pi \sum_0^n [\phi_{\nu}(0)]^* \phi_{\nu}(w)$$

so that  $2\pi \sum_0^n [\phi_{\nu}(0)]^* \phi_{\nu}(w)$  represents an analytic function in  $|w| < 1 + \epsilon$ . However it coincides with

$$\frac{1}{[D(0)]^* D(w)}$$

when  $|w| < 1$  (see [6]). We can then extend  $1/D(w)$  analytically into  $|w| < 1 + \epsilon$ . But then  $D(w)$  is analytic and different from zero in  $|w| < 1 + \epsilon$ . But we have except on a set of measure zero

$$f(\lambda) = |D(e^{i\lambda})|^2 = D(w)D^*(1/w), \quad w = e^{i\lambda},$$

where  $D^*(w)$  denotes the function obtained from  $D(w)$  by taking the conjugates of its Taylor coefficients. From this the result follows.

Let us note that Theorem 4 is true more generally for conditions of the type  $C_i$ .

We have seen that if  $\delta_n$  decreases exponentially,  $f(\lambda)$  cannot have any essential zeros. In the following theorem we study what happens when  $f(\lambda)$  has zeros.

**THEOREM 5.** *If  $f(\lambda)$  coincides almost everywhere with*

$$g(\lambda) \prod_{\nu=1}^s |e^{i\lambda} - e^{i\lambda_\nu}|^{2l_\nu},$$

where  $g(\lambda)$  is positive and has an integrable third derivative, then

$$\delta_n = O(1/n).$$

The order is attained for some such  $f(\lambda)$ .

**Proof.** If  $g(\lambda)$  has an integrable third derivative,  $D(w)$  has a bounded derivative on  $|w| = 1$ . Repeating an argument used in the proof of Theorem 2, we see that

$$\mu_n = \beta' (H'_{n+M})^{-1} \beta, \quad M = \sum l_\nu,$$

where  $\beta$  has its first component equal to one and the remaining components are zero and

$$H_{n+M} = \mathcal{D}_n \left\{ \begin{pmatrix} \Lambda & & & \\ & \gamma_1 & & 0 \\ & 0 & \ddots & \\ & & & \gamma_p \end{pmatrix} + \frac{1}{n} B_n \right\} \mathcal{D}_n$$

where  $B_n$  is bounded and

$$\mathcal{D}_n = \begin{pmatrix} 1 & & & & \\ & n^{1/2} & & & 0 \\ & & \ddots & & \\ & & & n^{h+1/2} & \\ 0 & & & & n^{1/2} \\ & & & & & \ddots \\ & & & & & & \ddots \end{pmatrix}.$$

The theorem follows immediately. The bound can be realized by  $f(\lambda) = |1 - e^{i\lambda}|^2$ .

**5. Applications to stochastic processes.** Consider a discrete stochastic process  $x_t, -\infty < t < \infty$ , with mean value  $m_t = Ex_t$ . We assume that the second order moments  $E|x_t|^2$  exist and that the reduced process  $y_t = x_t - m_t$

is stationary in the wide sense, that is

$$r_{s,t} = E y_s y_t^* = r_{s-t}.$$

Then we know that

$$r_\nu = \int_{-\pi}^{\pi} e^{i\nu\lambda} dF(\lambda)$$

where  $F(\lambda)$  is bounded and nondecreasing in  $(-\pi, \pi)$ . The completely non-deterministic processes form an important subclass. They are completely characterized by

(1)  $F(\lambda)$  absolutely continuous,

$$F(\lambda) = \int_{-\pi}^{\lambda} f(l) dl,$$

and

$$(2) \quad \int_{-\pi}^{\pi} \log f(\lambda) d\lambda > -\infty.$$

The nonnegative function  $f(\lambda)$  is called the spectral density of the process.

In various problems one is interested in minimizing the variance of a linear form  $\sum_0^n c_\nu x_\nu$ , subject to some conditions on the  $c_\nu$ 's. But this variance is

$$\sum_{\nu, \mu=0}^n c_\nu c_\mu^* r_{\nu-\mu}$$

which is of the form (1.1). We shall consider some problems of this type.

1. Let us first assume  $m_i \equiv 0$ . Having observed  $x_1, x_2, \dots, x_n$  we want to form a linear combination  $\sum_1^n c_\nu x_\nu$  such that

$$E \left| x_0 - \sum_1^n c_\nu x_\nu \right|^2 = \min.$$

This is a familiar problem of extrapolation. This is the type of problem treated in the previous sections since we can write it as

$$\mu_n = \int_{-\pi}^{\pi} |P_n(e^{i\lambda})|^2 f(\lambda) d\lambda = \min,$$

$$P_n(0) = 1.$$

The only restraint is inside the unit circle and we then know that  $\mu_n$  tends to a positive value  $\mu$  as  $n \rightarrow \infty$ . From Theorem 1 we get

$$(5.1) \quad \mu = 2\pi |D(0)|^2 = 2\pi \exp \left\{ \frac{1}{2\pi} \int_{-\pi}^{\pi} \log f(\lambda) d\lambda \right\}.$$

The "predictor" is

$$\sum_{\nu=1}^n c_\nu x_\nu = \int_{-\pi}^{\pi} \sum_1^n c_\nu e^{i\nu\lambda} dZ(\lambda)$$

where  $Z(\lambda)$  is the orthogonal process corresponding to the stationary process  $x_t$  (see [2]). The limit in the mean of this stochastic variable as  $n \rightarrow \infty$  is

$$(5.2) \quad \int_{-\pi}^{\pi} \left[ 1 - \frac{D(0)}{D(e^{i\lambda})} \right] dZ(\lambda).$$

(5.1) and (5.2) can be found in [7].

2. A slightly more general case is extrapolation  $k$  steps back, i.e., given a sample  $x_k, x_{k+1}, \dots, x_n$  to predict  $x_0$ . We see that this corresponds to the conditions

$$\begin{aligned} P(0) &= 1, \\ P'(0) &= 0, \\ &\dots \dots \dots, \\ P^{(k-1)}(0) &= 0. \end{aligned}$$

This is again a context treated in Theorem 1. Analogues of expressions (5.1) and (5.2) in this case can be found in a similar way.

3. Suppose that  $m_t$  is equal to an unknown constant  $m$ . From the sample  $x_0, x_1, \dots, x_n$  we want to construct a linear, unbiased estimate  $m^*$  of minimum variance. It is immediately seen that this is equivalent to the minimization of  $\int_{-\pi}^{\pi} |P_n(e^{i\lambda})|^2 f(\lambda) d\lambda$  under the condition  $P_n(1) = 1$ . As this is a condition on the unit circle, we can apply Theorem 2 which gives us

$$E | m^* - m |^2 \sim \frac{2\pi f(0)}{n}.$$

Theorem 3 implies that if  $f(0) \neq 0$  and  $f(\lambda)$  is continuous, we get an asymptotically equivalent estimate by solving the same minimization problem for a uniform spectral density. But that would give us just the empirical mean

$$x^* = \frac{1}{n+1} \sum_0^n x_\nu.$$

4. Let  $m_t = mt(t-1) \dots (t-k+1)e^{it\lambda_0}$ . If we now want to get a linear unbiased estimate of minimum variance  $m^*$  of  $m$  we have the condition  $P_n^{(k)}(e^{i\lambda_0}) = e^{-ik\lambda_0}$ .

Again we can use Theorem 2 and get

$$E | m^* - m |^2 \sim \frac{2\pi f(e^{i\lambda_0})}{n^{2k+1}} (2k+1).$$

We can also apply Theorem 3.

5. If we are interested in polynomial or trigonometric regression we put

$$\phi_t^{(k,j)} = t(t-1) \cdots (t-k+1)e^{i\lambda_j t},$$

$$m_t = \sum_{k,j} c_{k,j} \phi_t^{(k,j)}$$

where the regression coefficients  $c_{k,j}$  are unknown. To get an unbiased minimum variance estimate of  $c_{k,j}$  we have the conditions

$$P^{(l)}(e^{i\lambda_s}) = 0 \quad \text{if } (l, s) \neq (k, j),$$

$$P^{(k)}(e^{i\lambda_j}) = e^{-ik\lambda_j}.$$

All the conditions are on the unit circle.

6. If  $m_t \equiv m$  is unknown and we wish to predict the value of  $x_0$  from  $x_k, x_{k+1}, \dots, x_n$  it may be advantageous to use an unbiased predictor

$$\sum_{v=k}^n c_v x_v, \quad \sum_{v=k}^n c_v = 1.$$

We then get the conditions

$$P(0) = 1,$$

$$P'(0) = 0,$$

$$\dots \dots \dots,$$

$$P^{(k-1)}(0) = 0,$$

$$P(1) = 0.$$

It follows from Remark 2 that the limiting variance of this predictor is the same as that of the predictor in problem 2.

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UNIVERSITY OF CHICAGO,  
CHICAGO, ILL.