Stationary Processes, 1963-64 K. Ito

#### VI. PREDICTION AND MOVING AVERAGE REPRESENTATION

Linear Problem. (discrete time parameter). Let  $x_n \in \mathbb{Z}$ , be a complex-valued weakly stationary stochastic sequence with mean 0. We have the following decompositions:

$$r(n) = r_x(n) = \int_{\Gamma} e^{i2\pi \lambda n} dF(\lambda)$$
  $\Gamma = R|Z$ 

$$x_n = \int_{\Gamma} e^{i2\pi\lambda n} dM(\lambda), (M(\Lambda_1), M(\Lambda_2)) = F(\Lambda_1 \cap \Lambda_2)$$

 $x_n$  is called <u>trivial</u> if  $\Gamma(0) = 0$ , i.e., if  $x_n = 0$ .

## 1. Definitions.

 $L_{min}(\mathbf{x}) = \text{closed linear subspace of } L^2(\Omega, \mathcal{B}, P) \text{ spanned by } x_k, m \le k \le n.$ 

$$L_{n}(x) = L_{-\infty,n}(x)$$

$$L(x) = L_{\infty}(x) = L_{-\infty,\infty}(x) = \bigvee_{n} L_{n}(x)$$

$$L_{-\infty}(x) = \Lambda L_n(x)$$
 (= the space of remote past),  $m < n \implies L_m(x) \subset L_n(x)$ 

shift operator U : unitary operator U determined by U  $x_n = x_{n+1}$ , n  $\in \mathbb{Z}$ ,  $U^n L_m(x) = L_{m+n}(x)$ 

Definition 1.  $x_n$  is called (purely) non-deterministic if  $L_{\infty}(x) = 0$  and (purely) deterministic if  $L_{\infty}(x) = L(x)$ .

Definition 2.  $x_n^d = P_{L_{\infty}}(x) \cdot x_n$ ,  $n \in \mathbb{Z}$ , is a deterministic stationary sequence and is called the <u>deterministic</u> part of  $x_n$ .

Definition 3.  $x_n^i = P_{L(x) \bigcirc L_{-\infty}(x)} \cdot x_n$  is a non-deterministic stationary sequence and is called the non-deterministic part of  $x_n$ .

Definition 4.  $x_n = x_n^1 + x_n^d$  is called the Wold decomposition of  $x_n$ 

Corollary 1.  $L(x^d) = L_{\infty}(x)$ ,  $L(x^i) = L(x) \bigcirc L_{\infty}(x)$ .

Definition 5. Two stationary sequences  $x_n$  and  $y_n$  which may be defined on different probability spaces are called (weakly) equivalent (in symbol  $x_n \sim y_n$ ) if  $r_x(n) \equiv r_y(n)$ .

Corollary 2.  $x_n \sim y_n$  iff  $\exists$  isomorphism  $\forall : L(x) \rightarrow L(y)$  with  $\forall x_n = y_n$ ,  $n \in \mathbb{Z}$ .

Corollary 3. A stationary sequence equivalent to a deterministic one is also deterministic. Similarly for "non-deterministic".

Definition 6. An orthonormal sequence  $\xi_n$  ( $(\xi_n, \lambda) = 0$ ,  $(\xi_n, \xi_m) = \delta_{nm}$ ) is called white light or white noise.

The Hincin measure of  $\xi_n$  is the uniform distribution on  $\Gamma$ :

$$r_{\xi}(n) = \delta_{no} = \int_{\Gamma} e^{i2\pi\lambda n} d\lambda,$$

from which the adjective "white" comes.

$$\frac{\text{Corollary 4}.}{\Gamma} = \int_{\Gamma} e^{i2\pi\lambda n} M_{\xi}(d\lambda), \quad (M(\Lambda_{1}), M(\Lambda_{2})) = \int_{\Lambda_{1}} \Lambda_{2} d\lambda$$

Corollary 5. A white noise is non-deterministic.

Corollary 6. A sequence equivalent to a white noise is also a white noise, and any two white noises are equivalent.

<u>Definition 7.</u> A stationary sequence  $x_n$  is said to have <u>linear regression</u> if it satisfies a linear difference equation with constant coefficients:

$$x_n + a_1 x_n + \cdots + a_m x_{n-m} = 0.$$

## Corollary 7.

x, has linear regression

<==> the Hincin measure increases only with a finite number of jumps.

# 2. Moving Average Representation.

Given a white noise  $\xi_n$  and a (non-random) two-sided sequence a  $E(a_n) \in L^2(Z)$ , form

$$y_n = \sum_m a_{n-m} \xi_m$$

i.e.,

Then y is also a stationary sequence with

$$r_y(n) = \int_{\Gamma} e^{i2\pi\lambda n} |\alpha(e^{-i2\pi\lambda})|^2 d\lambda$$

where

$$\mathbf{a}(\zeta) = \sum_{\mathbf{n} \in \mathbf{Z}} \mathbf{a_n} \zeta^{\mathbf{n}}$$

In fact

$$y_{n} = \sum_{m} a_{n \le m} \int_{\Gamma} e^{i2\pi\lambda m} M_{\xi}(d\lambda) = \int_{\Gamma} e^{i2\pi\lambda n} a(e^{-i2\lambda}) M_{\xi}(d\lambda)$$

so that

$$r_y(n) = (y_n, y_0) = \int_{\Gamma} e^{i2\pi\lambda n} |a(e^{-i2\pi\lambda})|^2 d\lambda$$

Theorem 1.  $L(a * \xi) = L(\xi)$  and  $L(a * \xi) = L(\xi) \iff e(e^{-i2\pi\lambda}) \neq 0$  a.e.

Proof. Use the following Tallberian Theorem for Fourier series: Let  $a \in \ell^2(Z)$ .

c. 
$$\ell$$
,  $m[\alpha * a, \alpha \in \ell^1(Z)] = \ell^2(Z)$ 
 $\iff a(e^{-i2\pi\lambda}) \neq 0, a.e.$ 

Definition 1. Let  $x_n$  be a stationary sequence. If we have

$$x \sim a * \xi$$
  $a \in \ell^2(Z)$ ,  $\xi = \text{white noise}$ ,

a \* & is called a moving average representation of x.

Corollary 1. If x has a moving average representation a \* § with  $a(e^{-i2\pi\lambda}) \neq 0$ , a.e., then we can find a white noise  $\eta_n \in L(x)$  such that

$$x = a * \eta$$

Proof. If  $a(e^{-12\pi\lambda}) \neq 0$ , then

$$L(y) = L(\xi),$$
 for  $y = a * \xi$ 

Since  $x \sim y$ , we have an isomorphism  $V : L(x) \to L(y)$  such that  $Vx_n = y_n$ . Set  $\eta_n = V^{-1} \xi_n$ . Then  $y = a * \xi$  goes over into  $x = a * \eta$  by  $V^{-1}$ . Theorem 2. In order for  $x_n$  to have a moving average representation, it is necessary and sufficient that the Hincin measure of x is absolutely continuous.

Proof. (i) Assume that  $x \sim y = a + \xi$ . Then

$$r_x(n) = r_y(n) = \int_{\Gamma} e^{i2\pi\lambda n} |a(e^{-i2\pi\lambda})|^2 d\lambda$$

(ii) Assume that

$$(\mathbf{x}_{\mathbf{p}} \ \mathbf{x}_{2}) = \int_{\Gamma} e^{i2\pi(\mathbf{p}-2)\lambda} f(\lambda) d\lambda, \qquad f \in L^{1}(\Gamma)$$

Then  $\sqrt{f(\lambda)} \in L^2(\Gamma)$ . Consider the Fourier expansion of  $\sqrt{f(\lambda)}$ 

$$\sqrt{f(\lambda)} = \sum_{m} a_{m} e^{-12\pi\lambda m} = a(e^{-12\pi\lambda})$$

with  $a = (a_n) \in \ell^2(Z)$ . Let  $\xi_n$  be a white noise. Then

because

$$r_{y}(n) = \int e^{i2\pi\lambda n} |a(e^{-i2\pi\lambda})|^{2} d\lambda$$
$$= \int e^{i2\pi\lambda n} f(\lambda) d\lambda = r_{x}(n).$$

Remark 1. The representation is not unique. We can use the expansion of any function  $\sqrt{f(\lambda)} e^{i\phi(\lambda)}$ ,  $\phi(\lambda)$  being real and measurable.

Remark 2. 
$$f(\lambda) \neq 0$$
 (a.e.)  $\iff a(e^{-2\pi\lambda}) \neq 0$  (a.e.)  $\iff L(a * \xi) = L(\xi)$   $\iff x = a * \eta, \qquad \eta_n \in L(x), \ (\eta_n) = \text{ white noise.}$ 

Definition 2. A moving average representation  $x \sim a * \xi$  is called backward if  $a_n = 0$  for n < 0 and forward if  $a_n = 0$  for n > 0.

Corollary 2. If  $x \sim a * \xi$  is backward, then

$$L_n(a * \xi) \sum_{n} L_n(\xi)$$

Theorem 3. The following three conditions are equivalent for a non-trivial stationary sequence  $x_n$ .

- (i)  $\mathbf{x}_{n}$  has a backward moving average representation,
- (ii) x<sub>n</sub> is non-deterministic,
- (iii) the Rinčin measure of  $\mathbf{x}_n$  is absolutely continuous with the density  $f(\lambda)$  satisfying

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

Proof. (i)  $\Longrightarrow$  (ii)

$$\mathbf{L}_{\mathbf{n}}(\mathbf{a} + \mathbf{\xi}) \subset \mathbf{L}_{\mathbf{n}}(\mathbf{\xi})$$

and so

$$L_{\infty}(\mathbf{a} + \boldsymbol{\xi}) \subset L_{\infty}(\boldsymbol{\xi}) = 0$$

Therefore a \* & is non-deterministic and so is x.

 $(ii) \rightarrow (i)$ 

$$L_{-1}(x) \subset L_0(x)$$

$$\begin{bmatrix} \therefore L_{-1}(x) = L_{0}(x) \implies L_{n-1}(x) = U^{n} \ L_{-1}(x) = U^{n} \ L_{0}(x) \\ = L_{n}(x) \implies L_{\infty}(x) = L_{\infty}(x) = 0 \\ \implies x_{n} \text{ is trivial (contrary to the assumption)} \end{bmatrix}$$

... 
$$x_0 \notin L_{-1}(x)$$
,  $\xi_0' = x_0 - P_{L_{-1}}(x) \cdot x_0 \neq 0$ 

Set 
$$\xi_0 = \frac{\xi_0'}{\|\xi_0'\|}$$
 and  $\xi_n = u^n \xi_0$ . Then

$$L_0(x) = L_{-1}(x) \oplus \{\xi_0\},$$

and applying U to both sides, we get

$$L_{n}(x) = L_{n-1}(x) \bigoplus \{\xi_{n}\}$$

$$L_{0}(x) = L_{-\infty}(x) + c.l.m. \{\xi_{0}, \xi_{1}, \xi_{-2}, ...\}$$

$$= c.l.m. \{\xi_{0}, \xi_{-1}, \xi_{-2}, ...\}$$

$$x_0 = \sum_{m} a_m \xi_{-m}$$

where

$$\mathbf{a}_{\mathbf{m}} = \begin{cases} (\mathbf{x}_{0}, \ \boldsymbol{\xi}_{-\mathbf{m}}) & \mathbf{m} \ge 0 \\ 0 & \mathbf{m} < 0 \end{cases}$$

Then

$$x_{n} = U^{n} x_{0} = \sum_{m} a_{m} U^{n} \xi_{-m} = \sum_{m} a_{m} \xi_{n-m}$$
$$= \sum_{m} a_{n-m} \xi_{m}$$

To prove (i) <=> (iii) we shall use the following known

Theorem: In order for a non-negative f  $\in L^1(\Gamma)$  to be expressed as

$$f(\lambda) = \left| \sum_{n > 0} a_n e^{-i2\pi \lambda n} \right|^2$$

with  $0 < \frac{\pi}{n} |a_n|^2 < \infty$ ,  $a_n = 0$  for n < 0, it is necessary and sufficient that

(This theorem is (1) po to particular induction in the results in the Section 4)

(1) → (111) If x ~ a \* § (backward), then

$$r_{x}(x) = r_{a + \xi}(n) = \int_{\Gamma} e^{i2\pi \lambda n} \Big|_{n \ge 0} \sum_{n \ge 0} a_{n} e^{-i2\pi \lambda n} \Big|^{2} d\lambda$$

and so the Hinčin measure of  $\mathbf{x}_n$  is absolutely continuous with the density

$$f(\lambda) = \left| \sum_{n \geq 0} a_n e^{-i2\pi \lambda n} \right|^2$$

low use the above theorem.

(iii)  $\to$  (i) If  $\int_\Gamma \log f(\lambda) \ d\lambda > -\infty$ , then the above theorem shows that  $f(\lambda)$  can be expressed as

$$f(\lambda) = \left| \sum_{n} a_n e^{-i2\pi \lambda n} \right|^2 \qquad (a_n = 0 \text{ for } n < 0),$$

and so  $x \sim a + \xi$  ( $\xi = \text{white noise}^{\top}$  because

$$r_{n \neq \xi}(n) = \int e^{i2\pi\lambda n} f(\lambda) d\lambda = r_{\chi}(n).$$

#### 3. Generalized Poisson Formula.

Theorem of Fatou. Let u(z) be a harmonic function and suppose that

(1) 
$$\sup_{0 \le r < 1 - \pi} \int_{-\pi}^{\pi} |u(re^{i\theta})| d\theta < \infty.$$

Then u can be expressed by generalized Poisson formula,

(2) 
$$u(re^{i\theta}) = \frac{1}{2\pi} \int_{-\pi}^{\pi} P_{r}(\theta - \phi) \mu(d\phi)$$

$$P_{r}(\mu) = \frac{1 - r^{2}}{1 - 2r \cos \theta + r^{2}} \qquad (Poisson kernel)$$

$$\mu = \psi - \lim_{r \to 1} \mu_{r}, \qquad \mu_{r}(d\theta) = u(re^{i\theta}) d\theta$$

The boundary value  $u(e^{i\theta})$  of u(z) exists a.e., and equals the density of the absolutely continuous part of  $\mu_{\theta}$  a.e. (To be more precise, where  $\mu'(\theta) = \lim_{\epsilon \to 0} \mu(\theta - \epsilon_1 \mu + \epsilon)/2\epsilon$  exists, the nontangential limit of u(z) as  $z \to e^{i\theta}$  exists and equals  $\mu'(\theta)$ .)

Note. In the general Poisson formula the a.e. existing boundary value  $u(e^{i\theta})$  determines only the absolutely continuous part of  $\mu(d\theta)$  and so does not always determine the behavior of u(z) in |z| < 1. To determine u(z), |z| < 1, completely, we should know, besides the boundary value  $u(e^{i\theta})$ , the singular part  $s(d\theta)$  of  $\mu(d\theta)$ , i.e., the weak limit of  $s_r(d\theta) = [u(re^{i\theta}) - u(e^{i\theta})] d\theta$  as  $r \neq 1$ .

Function Class  $H_1$  . (Hardy class of order 1)

Definition 1. A function analytic in |z| < 1 is said to belong to  $H_1$ , iff

(3) 
$$\sup_{0 \le r \le 1 - \pi} \int_{-\pi}^{\pi} |g(re^{i\theta})| d\theta < \infty$$

Applying Faton's theorem to the real and imaginary parts, we have

Generalized Poisson formula for  $H_1$ -functions. Any  $g \in H_1$  has boundary values  $g(e^{i\theta})$  a.e. and

(4) 
$$g(re^{i\theta}) = \frac{1}{2\pi} \int_{-\pi}^{\pi} Pr(\theta - \phi) [g(e^{i\phi}) d\phi + s(d\phi)]$$

where s is a complex-valued singular measure of bounded variation defined by

(5) 
$$s = w^* - \lim_{r \to 1} s_r, \qquad s_r(d\theta) = (g(re^{i\theta}) - g(e^{i\theta})) d\theta$$

Fundtion Class H (Hardy class of order 2).

Definition 2. A function analytic in |z| < 1 is said to belong to  $H_1$  iff

(6) 
$$\sup_{0 < r < 1 - \pi} |g(re^{i\theta})|^2 d\theta < \infty$$

Consider the power series expansion of g(z):

(7) 
$$g(z) = \sum_{n \geq 0} a_n z^n$$

Then (6) is equivalent to

(8) 
$$\sum_{n \geq 0} |a_n|^2 < \infty$$

Since  $g \in H_2 \subset H_1$ ,  $g(e^{i\theta}) = a.e. \lim_{r \to 1} g(re^{i\theta})$  exists. It follows from (7) that

(9) 
$$g(e^{i\theta}) = 1.i.m. \sum_{n \to \infty}^{N} a_n e^{in\theta} = 1.i.m.g(re^{i\theta})$$

from which we have

(10) 
$$\frac{1}{2\pi} \int_{-\pi}^{\pi} g(e^{i\theta}) e^{-in\theta} d\theta = \begin{cases} a_n & n \geq 0 \\ 0 & n < 0 \end{cases}$$

(11) 
$$g(re^{i\theta}) = \sum_{n \geq 0} a_n r^n e^{in\theta}$$

$$= \frac{1}{2\pi} \sum_{n=-\infty}^{\infty} r^{|n|} e^{in\theta} \int_{-\pi}^{\pi} g(e^{i\phi}) e^{-in\phi} d\phi$$

$$= \frac{1}{2\pi} \int_{-\pi}^{\pi} g(e^{i\phi}) \sum_{n=-\infty}^{\infty} r^{|n|} e^{in(\theta-\phi)} d\phi$$

and so

Then

(12) 
$$g(re^{i\theta}) = \frac{1}{2\pi} \int_{-\pi}^{\pi} P_r(\theta-\phi) g(e^{i\phi}) d\phi$$

This is Poisson formula for H, functions.

Relation Between H2 and L2.

Let  $L_{+}^{2}$  be the subspace of  $L^{2}(-\pi, \pi)$  generated by  $e^{in\theta}$ ,  $n \ge 0$ .

$$T: g(z) = \sum_{n \geq 0} a_n z^n \rightarrow (Tg) (e^{i\theta}) = \sum_{n \geq 0} a_n e^{in\theta} ,$$

$$\left(\sum_{n \geq 0} |a_n|^2 < \infty\right)$$

determines a one-to-one mapping from  $H_2$  onto  $L_+^2$  and we have

- (i) (Tg)  $(e^{i\theta})$  is the a.e. boundary value function of g and so can be written as  $g(e^{i\theta})$ .
- (11)  $\mathbf{g}_{\mathbf{r}}(e^{i\theta}) = \mathbf{g}(re^{i\theta}) \in \mathbf{L}_{+}^{2}$  and  $\|\mathbf{g}_{\mathbf{r}}\| + \|\mathbf{T}\mathbf{g}\|$  as r + 1. (Notice that  $\|\mathbf{g}_{\mathbf{r}}\| = \sum_{n} |\mathbf{a}_{n}|^{2} r^{2n}$ ).
- (111) Poisson Formula.

$$g(z) = \frac{1}{2\pi} \int_{-\pi}^{\pi} (Tg) (e^{i\Phi}) P_{r}(\theta - \phi) d\phi, \qquad z = re^{i\theta} \qquad (0 \le r < 1)$$

(iv) 
$$|g(z)| \le ||Tg|| \frac{1+|g|}{1-|g|}$$

(v) 
$$g(z) = \frac{1}{2\pi} \int_{-\pi}^{\pi} \P[(Tg)(e^{i\phi})] \cdot \frac{e^{i\phi} + z}{e^{i\phi} - z} d\phi + i\phi$$

We have already proved (i), (ii) and (iii).

To prove (iv),

$$\begin{split} |g(re^{i\theta})| &\leq \frac{1}{2\pi} \int_{-\pi}^{\pi} |Tg(e^{i\phi})| \frac{1+r}{1-r} \ d\phi \\ &\leq ||Tg|| \frac{1+r}{1-r} \end{split} \qquad \text{(by the Schwarz inequality)}$$

To prove (v), let h(z) be the integral in the right side.

Then

$$R_{h(z)} = \frac{1}{2\pi} \int_{-\pi}^{\pi} R[(2g) (e^{i\phi})] P_{r}(\theta - \phi) d\phi, \qquad z = re^{i\theta}$$

$$= Rg(z)$$

that which implies/h(z) - g(z) is a pure imaginary constant because h(z) and g(z) are analytic.

# 4. Factorisation theorem for functions in H2.

(1) 
$$g(s)=\alpha \cdot \prod_{k=1}^{n} \frac{z-\alpha_k}{1-\alpha_k z} \exp\left\{\frac{1}{2\pi}\int_{-\pi}^{\pi} \frac{e^{i\varphi}+z}{e^{i\varphi}-z} \log |g(e^{i\varphi})| d\varphi\right\},$$

where of is a constant of modulus 1, and

$$(2) \frac{1}{2\pi} \int_{-\pi}^{\pi} \log |\mathbf{g}(\mathbf{e}^{\mathbf{i}\uparrow})| d\varphi = \log |\mathbf{e}| + \sum_{\substack{k=1 \ \alpha_{k} \neq 0}}^{n} \log \frac{1}{|\mathbf{x}|},$$

where a is the first non-vanishing coefficient of the power series expansion of g(z).

The purpose of this section is to extend this theorem to functions in  $\mathbf{H}_2$ .

Lemma 1. Let g(z) be analytic in |z| < 1 and assume (that  $g(z) \neq 0$  or |z| < 1) and that  $g(0) \neq 0$ . Then if  $\alpha_1, \ldots, \alpha_n$  are the zero points of g in |z| < r,

$$\frac{1}{2\pi}\int_{-\pi}^{\pi}\log\left|\mathbf{g}(\mathbf{r}e^{i\varphi})\right|d\varphi = \log\left|\mathbf{g}(0)\right| + \sum_{\substack{k=1 \ \alpha_{k} \neq 0}}^{\pi}\log\frac{\mathbf{r}}{|\alpha_{k}|} \geq \log\left|\mathbf{g}(0)\right|.$$

Unless otherwise stated, we repeat every multiple root by its multiplicity.

2. If  $g \in H_2$  and if  $g(0) \neq 0$ , then

(4) 
$$\int_{-\pi}^{\pi} |\log|g(re^{i\varphi})| d\varphi$$
  
 $\leq \int_{-\pi}^{\pi} |g(e^{i\varphi})|^2 d\varphi - 2\pi \log|g(0)|$  for  $0 \leq r \leq 1$ .

3. If  $g \in H_2$ , then

$$|g(re^{i\theta})|^{2} \leq \frac{1}{2\pi} \int_{\pi}^{\pi} |g(e^{i\theta})|^{2} P_{r}(\theta-\theta) d\theta$$

$$|g(re^{i\theta})|^{2} \int_{\pi}^{\pi} |g(e^{i\theta})|^{2} P_{r}(\theta-\theta) d\theta \leq \frac{1}{2\pi} \int_{\pi}^{\pi} |g(e^{i\theta})|^{2} P_{r}(\theta-\theta) d\theta$$

$$|g(re^{i\theta})|^{2} \leq \frac{1}{2\pi} \int_{\pi}^{\pi} |g(e^{i\theta})|^{2} P_{r}(\theta-\theta) d\theta$$

4. If  $g \in H_2$  and if  $\alpha_1, \alpha_2, \ldots$  be the zero points of in |x| < 1, then

$$(7) \qquad \sum_{n \geq 0} (1 - |\alpha_n|) < \infty.$$

Lemma 5. If 
$$|\alpha_n| < 1$$
 and if  $\sum_{n} (1 - |\alpha_n|) < \infty$ ,

then the infinite product

(8) 
$$\int_{n}^{-\overline{\alpha}_{n}} \frac{z - \alpha_{n}}{|\alpha_{n}|} \frac{z - \alpha_{n}}{1 - \overline{\alpha}_{n} z} \qquad (convention: \quad \frac{\overline{\alpha}_{n}}{|\alpha_{n}|} = 1 \text{ if } \alpha_{n} = 0)$$

defines an enalytic function  $B(s) \in H_2$  with the zero points  $\{\alpha'_n\}$  and satisfies

$$|B(z)| \le 1 \quad \text{in} \quad |z| < 1,$$

(10) 
$$|B(e^{i\phi})| = 1$$
 a.e.

Definition 1. B(z) in Lemma 5 is called the Blaschke product with tero points  $\{\alpha_n\}$ .

Given  $g \in H_2$ , let  $\{w_n\}$  be its zero points in |z| < 1. Then we can define the Blaschke product with the roots  $\{w_n\}$  by Lemma 3 and Lemma 4. Then we have

Lemma 6. h m g/B  $\in$  H<sub>2</sub> and has no sero points in |s|<1 and |h(e<sup>14</sup>)| = |g(e<sup>14</sup>)| s.e.

If  $h \in H_2$  and if h has no roots in  $|\mathbf{x}| < 1$ , then

Lemma 7.  $h(\mathbf{x}) = \alpha \cdot \exp\left\{\frac{1}{2\pi} \int_{-\pi}^{\pi} \frac{\mathbf{a}^{1\gamma} + \mathbf{x}}{\mathbf{a}^{1\gamma} - \mathbf{x}} \left[\log |h(\mathbf{a}^{1\gamma})| \, d\gamma - \sigma(\alpha\gamma)\right]\right\}$ ,

where  $\propto$  is a constant of modulus 1 and  $\sigma$  is a bounded non-negative singular measure defined by  $d\sigma = \frac{v^{*}-\lim \left(\log |\mathbf{h}(e^{i\phi})| - \log |\mathbf{h}(re^{i\phi})|\right) d\phi}{r+1}$ .

<u>Perfinition</u> 2. A function f(s) is called an <u>outer function</u> with the (real) generating density  $\omega(e^{i\varphi})$ , if it is expressed as

(12) 
$$f(z) = \exp\left\{\frac{1}{2\pi}\int_{-\pi}^{\pi} \frac{e^{i\varphi} + z}{e^{i\varphi} - z} \omega(e^{i\varphi})d\varphi\right\}, \quad \omega(e^{i\varphi}) \in L^{1}(-\pi, \pi).$$

Notice that an outer function does not always belong to  ${\rm H}_2$  .

<u>Definition</u> 5. A function f(s) is called a <u>singular function</u> with <u>generating singular measure</u> do, if it is expressed as

(15) 
$$f(z) = \exp\left\{\frac{1}{2\pi}\int_{-\pi}^{\pi} \frac{e^{i\varphi} + z}{e^{i\varphi} - z}(-\sigma)(d\varphi)\right\},$$

where  $\sigma$  is a bounded non-negative singular measure on  $(-\pi,\pi)$ . Notice that a singular function belongs to  $H_2$ .

Definition 4. A function f analytic in |z| < 1 is called an inner in |z| < 1 function if  $|f(z)| \le 1$  (from which it follows that  $f \in H_2$ ) and if  $|f(e^{i\varphi})| = 1$  a.e. and the second

Corollary 1. Both Blaschke products and singular functions are inner;

inclary 2 
$$\frac{9}{4}$$
  $\frac{1}{4}$   $\frac{1}$ 

Theorem 2 (Factorization theorem for functions in H2).

Any function  $g \in H_2$  can be factorised as

(14)  $g = \alpha \cdot B \cdot g_a \cdot s$ 

of a constant of modulus 1

B : Blaschke product

gm : an outer function in H2

s : a singular function

g in the form (14), then

g is the Blaschke product with the same zero points as g,

B is the Blaschke product with the same zero points as g,  $g_g$  is the outer function with generating density  $\log |g(e^{i\phi})|$ , and

s is the singular function with generating singular measure  $\sigma(d\phi)$  defined by

(15) 
$$d\sigma(\varphi) = \psi^*-\lim_{r \uparrow 1} \left[\log|g(e^{i\varphi})| - \log|g(re^{i\varphi})|\right] d\varphi$$
.

We have also

(16) 
$$\frac{1}{2\pi} \int_{-\pi}^{\pi} \log|g(e^{i\varphi})| d\varphi = \log|a| + \sum_{\alpha_{k} \pm 0} \log \frac{1}{|\alpha_{k}|} + \int_{-\pi}^{\pi} d\sigma(\varphi),$$

where a is the first non-vanishing coefficient of the power series

expansion of g.

g is an inner function with

Corollary 2. g ∈ H. is a Blaschke product iff log g (re<sup>14</sup>) uniform

Corollary 2.  $g \in H_2$  is a Blaschke product iff  $\log |g(re^{i\varphi})|$  uniformly, integrable on  $[-\pi,\pi)$ .

Corollary 3.  $g \in H_2$  is an outer function iff  $\log |g(re^{i\varphi})|$  uniformly (in  $0 \le r < 1$ ) integrable on  $[-\pi,\pi)$ . and g(z) has no properties |g(z)| = |g(z)| = |g(z)| = |g(z)| = |g(z)|. Corollary 4.  $g \in H_2$  is a singular function iff g is an inner function which has no zero points in |z| < 1.

# 5. Determining all backward moving average representations of any non-deterministic stationary sequence.

In this section we shall determine all backward representations of the form

(1) 
$$x = a * \xi$$
,  $L(x) = L(\xi)$ ,  $U \xi_n = \xi_{n+1}$ 

for any (non-trivial) non-deterministic stationary sequence x<sub>n</sub>.

If  $x_n$  has any backward representation, then  $x_n$  must be non-deterministic and  $x_n$  has a unique backward representation of the form (1) with the same coefficients. This will justify that we consider only non-deterministic stationary sequences and only backward representations of the form (1).

Since  $x_n$  is non-deterministic, its Hinčin measure  $dF(\lambda)$  is absolutely continuous with density function  $f(\lambda)$  satisfying

(2) 
$$\int_{\Gamma} \log f(\lambda) d\lambda > -\infty$$

or equivalently (because of the integrability of f)

$$\int_{\Gamma} |\log f(\lambda)| d\lambda < \infty$$

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In Section 2 we saw that  $x = a + \xi$  is a backward representation,

(3) 
$$\begin{cases} \mathbf{a_n} = 0 & n < 0 \\ \Sigma |\mathbf{a_n}|^2 < \infty \\ |\Sigma \mathbf{a_n}|^2 = f(\lambda) & \text{a.e.} \end{cases}$$

Using  $\{a_n\}$ , we shall introduce an analytic function

(4) 
$$\mathbf{a}(\mathbf{z}) = \sum_{\mathbf{n}} \mathbf{a}_{\mathbf{n}} \mathbf{z}^{\mathbf{n}} \equiv \sum_{\mathbf{n} > 0} \mathbf{a}_{\mathbf{n}} \mathbf{z}^{\mathbf{n}}$$

Equation (3) can be written in terms of a(z) as

(3') 
$$\begin{cases} a \in H_2 \\ |a(e^{-i2\pi\lambda})| = |\lim_{r \to 1} a(re^{-i2\pi\lambda})| = \sqrt{f(\lambda)} \quad \text{a.e.} \end{cases}$$

By the factorization theorem in Section 4, a can be expressed as

(5) 
$$a(z) = \alpha B(z) \cdot s(z) \cdot g_{s}(z)$$

where B is the Blaschke product, s is a singular function and  $g_s$  is an outer function with generating density  $\log \sqrt{f(\lambda)}$  ( $\equiv \log |a(e^{-i2\pi\lambda})|$ )

Write  $g_i$  for  $\alpha \cdot B(z) \cdot s(z)$ . Then  $g_i(z)$  is an inner function which will be called the inner part of a, while  $g_i$  will be called the (standard) outer part of a.

Theorem 1. If x is a white noise and if  $x = a * \xi$  is a backward representation, then a(z) is an inner function and  $\xi$  can be expressed by x as

(6) 
$$\xi = a^* * x$$
,  $a_n^* = \overline{a}_{-n}$ 

i.e.,

(7) 
$$\xi_{n} = \sum_{k>0} \overline{a}_{k} x_{n+k}$$
 (forward representation)

Proof. Let  $\xi_n = \int_{\Gamma} e^{\frac{1}{2\pi}\lambda n} dN(1)$  be the Kolmogorov-Cramer representation of  $\xi_n$ . Since  $\xi_n$  is a white noise,

$$(M(\Lambda_1), M(\Lambda_2)) = \int_{\Lambda_1} \int_{\Omega} \Lambda_2 d\lambda$$

Then

$$x_n = \int_{\Gamma} e^{i2\pi\lambda n} a(e^{-i2\pi\lambda}) dM(\lambda)$$

Since x is also a white noise, we have

$$|a(e^{-12\pi\lambda})|^2 d\lambda = d\lambda$$

i.e.,  $|a(e^{-i2\pi\lambda})|=1$  a.e. But  $a(z)\in H_2$ , and so a(z) is inner. Now observe

$$k \ge 0 \quad \overline{a}_k \quad x_{n+k} = \int_{\Gamma} e^{i2\pi \lambda n} \quad \overline{\lambda}_k = e^{i2\pi \lambda n} \int_{R} e^{i2\pi \lambda$$

which proves (7).

If  $x = a * \xi$  is a backward representation, then we have

(8) 
$$L_{n}(x) \subset L_{n}(\xi)$$

and

$$L(x) = L(\xi)$$

Equation (8) is clear because  $x = a * \xi$  is backward. To prove (9), notice that

$$\left|\sum_{n>0} a_n e^{-i2\pi \lambda n}\right|^2 = f(\lambda) \neq a.e.$$

by virtue of (2), so that (b \* a : b  $\in L^1(Z)$ ) is dense in  $L^2(Z)$  by the Tauberian theorem.

Keeping (8) and (9) in mind, we shall introduce

Definition 1. A backward representation  $x = a^{-x} \xi$  is called canonical if

(10) 
$$L_n(x) = L_n(\xi)$$
 for every  $n \mapsto (-L_n(x) = L_n(\xi))$ 

There exists at least one canonical representation of any given non-deterministic stationary sequence  $x_n$ . Consider the innovation

$$f_g^u = \frac{\|x^u - \mathbf{h}^{r+1}(\mathbf{x}) \cdot \mathbf{x}^u\|}{\mathbf{h}^{r-1}(\mathbf{x}) \cdot \mathbf{x}^u}$$

and set

$$\mathbf{a}_{n}^{\mathbf{S}} = (\mathbf{x}_{0}, \boldsymbol{\xi}_{-n}^{\mathbf{S}}).$$

Then it is easy to see that  $x = a^8 * \xi^8$  is a canonical backward representation.

Definition 2.  $x = a^{8} * \xi^{8}$  is called a standard backward representation.

Theorem 2.

(11) 
$$\mathbf{a}^{\mathbf{E}}(\mathbf{Z}) = \exp \left\{ \int_{\mathbf{r}} \frac{e^{-\mathbf{1}2\pi\lambda} + \mathbf{Z}}{e^{-\mathbf{1}2\pi\lambda} + \mathbf{Z}} \log \sqrt{f(\lambda)} \, d\lambda \right\}$$

Proof. Let

(12) 
$$\mathbf{a}^{\mathbf{a}}(\mathbf{z}) = \mathbf{a}^{\mathbf{i}}(\mathbf{z}) \ \mathbf{a}^{\mathbf{o}}(\mathbf{z})$$

be the factorization of  $a^6$  into its inner and outer parts. Since  $|a^8(e^{-i2\pi\lambda})|^2 = f(\lambda)$  a.e.,  $a^0(z)$  must equal the right side of (11).

To prove (11), it is enough to prove a (3) = 1. Set

(13) 
$$\xi = a^{\frac{1}{4}} + \xi^{\frac{1}{8}}$$

Using the Kolmogorov-Cramer representation of  $\xi^8$ ,  $\xi$  is also a white noise and

$$x = a^{\circ} * \xi$$

By Theorem 1 we can derive from (13)

(15) 
$$\xi_n^s = \prod_{k > 0} \overline{a_k^i} \, \xi_{n+k}$$

Since  $\xi^8$  is the innovation of x, we have

$$\xi_n^s \in L_n(x) \subset L_n(\xi)$$
(by (14))

i.e.,  $\xi_n^s \perp \xi_{n+k}$  (k > 0), which, combined with (15), implies

(16) 
$$a_{k}^{1} = 0$$
 (k > 0) i.e.,  $a^{1}(z) = a_{0}^{1}$ 

It is clear that  $|a_0^1| = 1$  ( ...  $\xi_0 = a_0^1 \xi^6$ ). Thus  $a^8(z) = a_0^1 \cdot a^0(z)$ . But it is clear that

$$a^{0}(0) > 0,$$
  $a^{8}(0) > 0$ 

and so  $a_0^1 = 1$ .

Theorem 3. If  $x = a + \xi$  is any backward representation, and if  $a(z) = a^{1}(z) a^{0}(z)$  is the factorization of a(z) into the inner and outer parts of a(z) ( $a^{0}(z) = a^{0}(z)$  by Theorem 2), then

(17) 
$$\xi_{n} = \sum_{k \geq 0} \overline{a_{k}^{1}} \xi_{n+k}^{s}$$

Proof. Since  $a^i$  is inner,  $\eta = a^i * \xi$  is a white noise and  $x = a^s * \eta$ , so that  $\eta = \xi^s$ , i.e.,  $\xi^s = a^i * \xi$ , which implies (17) by Theorem 1.

Corollary 2.  $x = a + \xi$  is a canonical backward representation, iff  $a(z) = \alpha a^{\beta}(z)$ ,  $\alpha$  being a constant of modulus 1.

Remark. Let  $x = a * \xi$  be any backward representation. Then  $L(x) = L(\xi)$  and  $L_n(x) \subset L_n(\xi)$ . Therefore  $\xi_n \in L(\xi) \odot L_{n-1}(\xi) \subset L(x) \odot L_{n-1}(x)$  =  $L(\xi_n^s, \xi_{n+1}^s, \dots)$ . Therefore  $\xi_n = \sum_{k \geq 0} b_k \xi_{n+k}^s$ . (b<sub>k</sub> does not depend on n because of (1).) Theorem 3 shows that we can express  $b_k$  as  $a_k^T$  using the inner part  $a^1(x)$  of a(x).

In the canonical backward representation x = a \* we do not need the future information of x to construct  $\xi$ .

6. Prediction. Let  $\mathbf{x}_{\mathbf{n}}$  be any stationary sequence with the Hinčin measure of and the Kolmogorov-Gramer orthogonal random measure  $\mathfrak{M}(\lambda)$ . Let  $\mathbf{R}^1 = \mathbf{A} - \mathbf{S}$ ,  $\mathbf{A} \cap \mathbf{S} = \emptyset$  be the decomposition such that  $\mathbf{F}_{\mathbf{a}}(\Lambda) = \mathbf{F}(\Lambda \cap \mathbf{A})$  is absolutely continuous and  $\mathbf{F}_{\mathbf{a}}(\Lambda) = \mathbf{F}(\Lambda \cap \mathbf{S})$  is singular. It is clear that  $\mathbf{F}'(\lambda) = \mathbf{F}'_{\mathbf{a}}(\lambda)$  (=  $\mathbf{f}(\lambda)$ ),  $\mathbf{F}'_{\mathbf{s}}(\lambda) = \mathbf{0}$  a.e. Set

(1) 
$$M_{g}(\Lambda) = M(\Lambda \cap A), \qquad M_{g}(\Lambda) = M(\Lambda \cap S)$$

and

(2) 
$$x_n^a = \int e^{-i2\pi\lambda n} dM_a(\lambda), \quad x_n^s = \int e^{-i2\pi\lambda n} dM_s(\lambda)$$

Then

$$x_{n} = x_{n}^{a} + x_{n}^{s}$$

$$L(x^{a}) \perp L(x^{s})$$

$$L(x) = L(x^{a}) \bigoplus L(x^{s})$$

Recall that  $x_n$  can be decomposed as

$$x_n = x_n^1 + x_n^d, x_n^1 : non-deterministic, x_n^d : deterministic$$

$$(4) L(x^1) \perp L(x^d)$$

$$L(x) = L(x^1) \bigoplus L(x^d)$$

Let us examine the relation between two decompositions (3) and (4).

Theorem 1. If  $\int \log f(\lambda) d\lambda = -\infty$ , then  $x_n^a$ ,  $x_n^a$  and  $x_n^a$  are all deterministic and

$$x_n^i = 0,$$
  $x_n^d = x_n = x_n^a + x_n^s$ 

If  $\int \log f(\lambda) d\lambda > -\infty$ , then

$$x_n^1 = x_n^a, \qquad x_n^d = x_n^s$$

Definition 1.  $x_{n,m} = P_{L_m} \cdot x_n$  is called the <u>predictor</u> of  $x_n$  with the information up to time m, and  $e_{n,m}^2 = \|x_n - P_{L_m}(x) \cdot x_n\|^2$  is called the mean square error of this predictor  $e_{n,m}^2$  depends only on n-m and so we can write  $e_{n-m}^2$  for  $e_{n,m}^2$ .

If  $\int f'\lambda$  d $\lambda = -\infty$ , then  $x_n$  is deterministic and so  $x_{n,m} = x_n$  for every m. Therefore there is no problem of prediction in this case.

Theorem 2. If  $\int f(\lambda) d\lambda > -\infty$ , then

$$e_m^2 = \sum_{k=0}^{m-1} |a_k|^2$$

where  $a_k$ ,  $k = 0, 1, 2, \dots$  are determined by

$$\sum_{k \geq 0} a_k z^k = \exp \left\{ \frac{1}{2} \int_{\Gamma} \frac{e^{-i2\pi\lambda} + z}{e^{-i2\pi\lambda} - z} \log f(\lambda) d\lambda \right\}$$

The predictor xn, m is obtained as follows:

$$\xi_n = \frac{x_n - P_{n-1}(x) \cdot x_n}{|x_n - P_{n-1}(x) \cdot x_n|} \quad (innovation)$$

$$a_{k} = (x_{0}, \xi_{-k}) = (x_{n}, \xi_{n-k})$$

$$y_{n,m} = \sum_{k \ge n-m} a_k f_{n-k} + x_n^s$$
(see the second)

# 7. Concrete expression of innovation and predictor.

Let  $x_n$  be non-deterministic. Then the Hinčin measure is absolutely continuous with the density  $f(\lambda)$  satisfying

(1) 
$$\int \log f(\lambda) d\lambda > -\infty$$

We shall introduce the following functions analytic in  $|\mathbf{z}| < 1$ 

(2) 
$$\mathbf{a}(\mathbf{z}) = \mathbf{E} \mathbf{a}_{\mathbf{k}} \mathbf{z}^{\mathbf{k}} = \exp \left\{ \frac{1}{2} \int_{\Gamma} \frac{e^{-i2\pi\lambda} + \mathbf{z}}{e^{-i2\pi\lambda} - \mathbf{z}} \log f(\lambda) d\lambda \right\}$$

(3) 
$$b(z) = \sum b_k z^k = a(z)^{-1}$$

(4) 
$$A_{\mathbf{z}}(z) = \sum_{k \geq \mathbf{z}} a_k z^k$$

(5) 
$$C_{\ell}(z) = \sum_{k > \ell} c_{\ell,k} z^{k} = A_{\ell}(z) b(z)$$

It is clear that a(z) is a standard outer function and belongs to  $H_2$ . Therefore  $\Sigma |a_n|^2 < \infty$ , so that  $A_{\ell}(z) \in H_2$ . b does not necessarily belong to  $H_2$ , but the a.e. boundary value function  $b(e^{-i2\pi\lambda})$  exists and equals  $a(e^{-i2\pi\lambda})^{-1}$  a.e. Similarly for  $C_{\ell}(z)$ .

Now observe the Kolmogorov-Cramer representation for  $\mathbf{x}_n$  and  $\mathbf{t}_n$ 

(6) 
$$x_n = \int e^{i2\pi \lambda n} dx_x(\lambda)$$

(7) 
$$\mathbf{k_n} = \int e^{\mathbf{i}2\pi\lambda \mathbf{n}} d\mathbf{M}_{\mathbf{g}}(\lambda)$$

Then it is clear that

(9) 
$$dM_{\xi}(\lambda) = b(e^{-i2\pi\lambda}) dM_{\chi}(\lambda)$$

Thus we have

(10) 
$$\xi_n = \int e^{i2\pi\lambda n} b(e^{-i2\pi\lambda}) dM(\lambda)$$

(11) 
$$x_{n,n-\ell} = \int e^{i2\pi\lambda n} A_{\ell} (e^{-i2\pi\lambda}) dM_{\ell}(\lambda)$$

$$= \int e^{i2\pi\lambda n} A_{\ell} (e^{-i2\pi\lambda}) b(e^{-i2\pi\lambda}) dM(\lambda)$$

$$= \int e^{i2\pi\lambda n} C_{\ell} (e^{-i2\pi\lambda}) dM(\lambda)$$

from which we have formal expansions

(10') 
$$\xi_{n} = \sum_{k > 0} b_{k} x_{n-k}$$

(11') 
$$x_{n,n-\ell} = \sum_{k \geq \ell} c_{\ell,k} x_{n-k}$$

Theorem 1. If  $f(\lambda)$  is essentially bounded (i.e.,  $\exists M < \infty$  such that  $f(\lambda) < M$  a.e.), and if  $\int f(\lambda)^{-1} d\lambda < \infty$ , then the formal expansions (10'), (11') converge and the equalities are true.

Proof. It follows from (2) and (3) that

$$b(z) = \exp\{-\frac{1}{2} \int \cdots \log f(\lambda) d\lambda\}$$
$$= \exp\{\frac{1}{2} \int \cdots \log f(\lambda)^{-1} d\lambda\}$$

and  $\int f(\lambda)^{-1} d\lambda < \infty$  implies that  $b \in H_2$ , so that  $\sum |b_p|^2 < \infty$  and  $b(e^{-i2\pi\lambda}) = \sum_{k>0} b_k e^{-i2\pi\lambda k}$  (convergence in  $L^2(\Gamma d\lambda)$ )

$$\begin{aligned} &\|\boldsymbol{\xi}_{n} - \sum_{k=0}^{p} b_{k} \mathbf{x}_{n-k}\|^{2} \\ &= \| \int e^{12\pi\lambda n} \sum_{k>p} b_{k} e^{-12\pi\lambda k} dM(\lambda) \|^{2} \\ &= \int \left| \sum_{k\geq p} \cdots \right|^{2} f(\lambda) d\lambda \\ &\leq M \int \left| \sum_{k>p} \cdots \right|^{2} d\lambda = M \sum_{k>p} \left| b_{k} \right|^{2} \to 0 \qquad (p \to \infty) \end{aligned}$$

Since we have

$$|b(e^{-i2\pi\lambda})|^2 = f(\lambda)^{-1}$$

and

$$C_{\underline{\beta}}(z) = 1 - b(z) \sum_{k < \underline{\beta}} a_k z^k$$

we have

$$\int |C_{\underline{z}}(re^{-i2\pi\lambda})|^2 d\lambda \le 2 \int \left[1 + f(\lambda)^{-1} \left(\sum_{k \le \underline{z}} |a_k|\right)^2\right] d\lambda$$

and therefore  $C_{\underline{z}}(z) \in \mathbb{H}_2$ . By the same argument as for  $\xi_n$  we have

$$\|\mathbf{x}_{\mathbf{n},\mathbf{n}-L} - \sum_{\mathbf{k}=L}^{\mathbf{p}} \mathbf{c}_{L,\mathbf{k}} \mathbf{x}_{\mathbf{n}-\mathbf{k}}\|^{2}$$

$$\leq \mathbf{M} \sum_{\mathbf{k} \geq \mathbf{p}} |\mathbf{c}_{L,\mathbf{k}}|^{2} \to 0 \qquad (\mathbf{p} \to \mathbf{w})$$

In some cases the formal expansion does not converge but its Césaro sum converges and equals the right value, for example

Theorem 2. If 
$$f(\lambda) = |1 - e^{-i2\pi\lambda}|^2 = 2(1 - \cos 2\pi\lambda)$$
, then 
$$a(z) = 1 - z \qquad (standard outer function)$$
$$x_n = \xi_n - \xi_{n-1} \qquad (standard representation)$$

w and

$$b(z) = \frac{1}{1-z} = 1 + z + z^2 + \cdots$$

The formal expansion for  $\xi_n$  i.e.,  $x_n + x_{n-1} + x_{n-2} + \cdots$  does not converge but we have, for its Césaro sum,

$$\lim_{q \to \infty} \frac{1}{q} \prod_{p=0}^{q-1} (x_n + x_{n-1} + \cdots + x_{n-p}) = \xi_n$$

(Notice that  $x_{n,n-1} = -\xi_{n-1} = -\lim_{q \to \infty} \frac{1}{q} \sum_{p=0}^{q-1} (x_{n-1} + \cdots + x_{n-p-1})$  and  $x_{n,n-4} = 0$  for 4 > 1.)

Proof.

$$\frac{1}{q} \sum_{p=0}^{q-1} (x_n + x_{n-1} + \cdots + x_{n-p}) = \xi_n - \frac{1}{q} (\xi_{n-1} + \xi_{n-2} + \cdots + x_{n-q})$$

and the norm of the second term is  $q^{-1/2}$ .

# 8. Linear difference equation.

Let us consider a linear difference equation

(1) 
$$a_0 x_n + a_1 x_{n-1} + \cdots + a_k x_{n-k} = \xi_n$$

where  $\xi_n$  is a given white noise.

Theorem 1. If there exists a stationary sequence  $x_n$  satisfying (1), then

(2) 
$$\mathbf{a}(\mathbf{z}) = \prod_{j=0}^{k} \mathbf{a}_j \mathbf{z}^j$$
 has no roots on  $|\mathbf{z}| = 1$ 

Conversely if (2) holds, then there exists a unique stationary sequence  $x_n$  satisfying (1).

Proof. Let us consider the Kolmogorov-Cramer representation of &

(3) 
$$\xi_{n} = \int e^{i2\pi \lambda n} dt_{\xi}(\lambda)$$

and assume that

(4) 
$$x_n = \int e^{i2\pi \lambda n} dM_x(\lambda)$$

satisfies (1). Then

(5) 
$$a(e^{-i2\pi\lambda}) dM_{\pi}(\lambda) = dM_{\pi}(\lambda)$$

and so

(6) 
$$|\mathbf{a}(e^{-12\pi\lambda})|^2 d\mathbf{r}(\lambda) = d\lambda$$
,  $d\mathbf{r}(\lambda) = the Hinčin measure of  $x_n$$ 

Let N be the set of zero points of  $a(e^{-i2\pi\lambda})$  (which is clearly a finite set) and let  $\tilde{a}G$  be the restriction of dF over  $N^C$ . Then

$$dF \geq dG = |a(e^{-12\pi\lambda})|^{-2} d\lambda.$$

If N  $\ni \lambda_0$ , then  $a(e^{-i2\pi\lambda}) \sim non-vanishing constant <math>\times (\lambda - \lambda_0)$  and so

$$\int dP \ge \int |a(e^{-i2\pi\lambda})|^{-2} d\lambda = -,$$

which is a contradiction. Thus N must be empty. This proves the first half of our theorem.

If (2) holds, then the above argument shows that

and

(8) 
$$b(e^{-i2\pi\lambda}) = a(e^{-i2\pi\lambda})^{-1} \in L^2(\Gamma, d\lambda)$$

so that x must be expressed as

(9) 
$$x_{n} = \int e^{i2\pi\lambda n} dM_{x}(\lambda) = \int e^{i2\pi\lambda n} b(e^{-i2\pi\lambda}) dM_{\xi}(\lambda)$$

and it is clear that  $x_n$ , thus defined, solves (1).

Theorem 2. Assume that (2) holds. Then  $b(e^{-i2\pi\lambda}) \in L^2(\Gamma, d\lambda)$  and it can be expanded in Fourier series

(10) 
$$b(e^{-12\pi\lambda}) = \prod_{j=-\infty}^{\infty} b_j e^{-12\pi\lambda j}$$

and the solution  $x_n$  of (1) is given by

$$x_{n} = \sum_{j} b_{j} \xi_{n-j}$$

If a(z) has no roots in |z| < 1 (and so in  $|z| \le 1$ ), then  $b_j = 0$  (j < 0) in (11) and (11) gives a canonical backward representation of  $x_n$ .

If a(z) has roots in |z| < 1, then (11) is not backward.

Proof. The first and third parts are clear. To prove the second part, observe that

$$b(z) = a(z)^{-1}$$

$$= a_0^{-1}(1 - \alpha_1 z)^{-1} (1 - \alpha_2 z)^{-1} \cdots (1 - \alpha_k z)^{-1}, (|\alpha_j| < 1)$$

is an outer function in H2.