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STUDIES IN PROBABILITY THEORY

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RESTRICTED COLLECTION

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CALCUTTA

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INTRODUCTION

This thesis consists of four chapters.

The impetus for the work in Chapter 1 comes from the concept of 'conditional atom' introduced by Neveu [19]. Here, using conditional atoms we generalize the concept of nonatomicity of measures. (We confine ourselves to probability measures). We obtain generalizations of results on nonatomic measures in [1], [3] and of Liapounoff's theorem.

The results in Chapter 2 have their origins in a paper by Boylan [7]. To study 'equiconvergence of martingales' Boylan introduced in [7] a metric on the space of complete sub \sigma-algebras of a probability space. A little later, Neveu showed in [19a], this metric space of complete sub \sigma-algebras is 'tight'; that is, if two sub \sigma-algebras are 'close' under this metric with one contained in the other, then there is a set (a conditional atom), with 'high' probability on which the traces of these \sigma-algebras coincide. In Chapter 2 we investigate what else this metric space is besides being 'tight'. We prove a host of results concerning the topological properties of this metric space; we also study an isomorphism problem.

Chapters 3 and 4 are devoted to problems in martingales. In Chapter 3 we prove a convergence theorem for 'fairer with time processes' (a generalization due to Blake of martingales

in [5]). In Chapter 4 we disprove with the help of an example a conjecture on singular martingales made by Luis Baez-Duarte in [14].

Some of the results in this thesis have already appeared in print. See [23] and [24].

CHAPTER 1

A GENERALIZATION OF NONATORICITY OF MEASURES

1. Introduction

Let $(\underline{\Gamma})$, \underline{A} , $\underline{\mu}$) be a probability space. Let $\underline{\mathbb{N}}_{\mu}$ denote the class of all u-null sets. In [19], Neveu gives the following definition that generalizes the classical concept of an atom of a probability space.

<u>Definition</u>: Let \underline{B} be a sub σ -algebra of \underline{A} . A set $\underline{A}_0 \in \underline{A}$ is called a conditional atom of \underline{A} with respect to \underline{B} if μ (\underline{A}_0) > 0 and the traces \underline{A}_0 (\underline{C}) $\underline{A} = \underline{A}_0$ (\underline{C}) \underline{B} (modulo $\underline{\underline{M}}_{\mu}$); (that is, if for every $\underline{A} \in \underline{A}$ one can find $\underline{B} \in \underline{B}$ such that $\mu(\underline{A}_0) = \underline{A}_0 = \underline{A}_0$ (\underline{C}) $\underline{A}_0 = \underline{A}_0 = \underline{A}_0$ (\underline{C}) $\underline{A}_0 = \underline{A}_0 = \underline{A}_0$ (\underline{C}) $\underline{A}_0 = \underline{A}_0 = \underline{A}_0 = \underline{A}_0$ (\underline{C}) atom.

It is easily verified that if $\underline{B} = \{\emptyset, \bigcap\}$, the conditional atom as defined above coincides with the classical concept of an atom. This generalization leads, in a natural way, to that of the concept of nonatomicity of a probability.

Definition: μ is said to be $(\underline{B}, \underline{A})$ nonatomic if there is no $(\underline{B}, \underline{A})(\mu)$ atom.

Clearly when $\underline{B} = \{ \emptyset, \widehat{\Box} \}$, this concept coincides with the classical concept of nonatomicity; so when μ is (\underline{B} , \underline{A})

nonatomic for $\underline{\mathbb{B}}=\{\varnothing,\ \bigcap\}$, we shall simply write μ is nonatomic on $\underline{\mathbb{A}}$.

In this chapter in section 2 we give some characterisations of (B, A) nonatomicity, in section 3 we consider the product problem and in section 4 we obtain a generalization of Liapounoff's theorem. Section 5 is devoted to a generalization of a result of Halmos. In section 6 we study this nonatomicity on Polish (complete separable metric) spaces.

In the sequel ($(\ \ \)$, $(\ \)$ will stand for a measurable space where $(\ \)$ is a σ -algebra. $(\ \)$, (

Following Boylan [7], we can introduce a pseudo-metric d_{μ} on \underline{S} , the space of all sub σ -algebras of the probability space ($\widehat{\Gamma}$), \underline{A} , μ), by setting

$$d_{\mu}(\underline{B}_{1}, \underline{B}_{2}) = \sup_{\substack{B_{1} \in B \\ 1 = 1}} \inf_{\substack{B_{2} \in B \\ 2 = 2}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{2} \in B \\ B_{2} \in B_{2}}} \inf_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{2} \in B \\ B_{2} \in B_{2}}} \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{1} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2}}} u(B_{1} \Delta B_{2}) + \sup_{\substack{B_{1} \in B_{2} \\ B_{2} \in B_{2$$

for \underline{B}_1 , $\underline{B}_2 \in \underline{S}$. This distance gives us an useful characterisation of nonatomicity. For our study we need the following result, a proof of which can be found in [19].

Theorem 1.1 Let μ be $(\underline{B}, \underline{A})$ nonatomic. Given any \underline{B} -measurable function f with $0 \le f \le 1$, we can find $\underline{A} \in \underline{A}$ such that $\underline{P}_{\mu}(\underline{A}|\underline{B}) = f$.

The following well-known results will be used without mention. Every countably generated σ -algebra is atomic. A measure on a countably generated σ -algebra is nonatomic if and only if the measure gives mass zero to every atom of the σ -algebra.

2. Some characterisations of monatomicity.

Theorem 2.1 Let μ be $(\underline{B}, \underline{A})$ nonatomic and $\lambda << \mu$. Then λ is also $(\underline{B}, \underline{A})$ nonatomic.

Proof: Let g denote a version of the Radon-Nikodym derivative of λ with respect to μ . Let $D = \{g > 0\}$. Observe that (i) λ (D) = 1 and (ii) for any $A \in \underline{A}$, $A \subset D$ we have λ (A) = 0 if and only if u(A) = 0. Now, let if possible A_0 be a (\underline{B} , \underline{A}) (λ) atom. By (i) λ ($A_0 \subset D$) = λ ($A_0 \subset D$) > 0. Also, since any subset of positive measure of a conditional atom is itself a conditional atom, it follows that $A_0 \subset D$ is a (\underline{B} , \underline{A}) (λ) atom.

That is, $(A_0 \cap D) \cap \underline{A} = (A_0 \cap D) \cap \underline{B} \pmod{\underline{N}}$.

This implies, in view of (ii), $(A_0 \cap D) \cap \underline{A} = (A_0 \cap D) \cap \underline{B} \pmod{\underline{N}}$. $(\text{modulo }\underline{\underline{N}}_{\mu})$. Again, by (ii), $\mu(A_0 \cap D) > 0$. So $A_0 \cap D$ is a $(\underline{B}, \underline{A}) \pmod{\mu}$ atom. A contradiction. Hence the theorem.

Before proceeding with the next result we would like to observe the following. Let $\underline{\mathbb{B}}_1 \subset \underline{\mathbb{B}}_2$. Then for any μ , $d_{\mu}(\underline{\mathbb{B}}_1, \underline{\mathbb{B}}_2) \leq 1/2$. This is because, for any $\underline{\mathbb{B}}_1 \in \underline{\mathbb{B}}_1$, $\inf_{B_1} \mu(\underline{\mathbb{B}}_1 \triangle \underline{\mathbb{B}}_2) = 0$; and for any $\underline{\mathbb{B}}_2 \in \underline{\mathbb{B}}_2$, $\inf_{B_1} \mu(\underline{\mathbb{B}}_1 \triangle \underline{\mathbb{B}}_2) \leq \underline{\mathbb{B}}_2$ $\lim_{B_1} \mu(\underline{\mathbb{B}}_2 \triangle \underline{\mathbb{B}}_2) \leq \underline{\mathbb{B}}_2$ $\lim_{B_1} \mu(\underline{\mathbb{B}}_2 \triangle \underline{\mathbb{B}}_2) \leq 1/2$.

Theorem 2.2 μ is (B, A) nonatomic if and only if for every $\lambda < \mu$, we have $d_{\lambda}(B, A) = 1/2$.

Proof: Let μ be $(\underline{B}, \underline{A})$ nonatomic. Let $\lambda << \mu$. Then by Theorem 2.1 λ is $(\underline{B}, \underline{A})$ nonatomic. By Theorem 1.1 there exists \underline{A} ε \underline{A} such that $\underline{P}_{\lambda}(\underline{A}, \underline{B}) = 1/2$. Now for any \underline{B} ε \underline{B} , λ $(\underline{A}, \Delta \underline{B}) = \lambda$ $(\underline{A}, -\underline{B}) + \lambda(\underline{B} - \underline{A})$

$$= \int_{\mathbb{B}^{c}} I_{A_{O}} d\lambda + \int_{\mathbb{B}} I_{A_{O}^{c}} d\lambda$$

$$= \int_{\mathbb{B}^{c}} P_{\lambda} \left(A_{O} \middle| B \right) d\lambda + \int_{\mathbb{B}} P_{\lambda} \left(A_{O}^{c} \middle| B \right) d\lambda$$

$$= (1/2) \lambda(B^{c}) + (1/2)\lambda(B) = 1/2$$

 $\inf_{B \in \underline{B}} \lambda \ (A \land B) = 1/2. \quad \text{So} \quad \frac{1}{2} \ge \sup_{A \in \underline{A}} \inf_{B \in \underline{B}} \lambda \ (A \land B) \ge 1/2.$

Hence $d_{\lambda}(B, A) = 1/2$ and this completes the 'only if' part.

Let A_0 be a $(\underline{B}, \underline{A})$ (u) atom. Define λ on \underline{A} by, for $\Lambda \in \underline{A}, \lambda(\Lambda) = \frac{\mu(\Lambda(\overline{A}))}{\mu(A_0)}$. Clearly $\lambda << \mu$. Since, for every $\Lambda \in \underline{A}$ there is a $B \in \underline{B}$ with $\mu[(\Lambda \wedge B)(\overline{A})] = 0$, we have $\lambda(\Lambda \wedge B) = 0$; hence $d_{\lambda}(\underline{B}, \underline{A}) = 0$ and this completes the 'if' part.

Remarks: With the same proof, therem 2.2 can be recast as $^{1}\mu$ is (B, A) nonatomic if and only if for every $\lambda << \mu$, we have d_{λ} (B, A) > 0'. The following stronger version of Theorem 2.2 can also be proved. μ is (B, A) nonatomic if and only if for every $\lambda \equiv \mu$ we have d_{λ} (B, A) = 1/2.

Theorem 2.3 Let \underline{B}_1 (\underline{B}_2 (\underline{B}_3) If μ is (\underline{B}_1 , \underline{B}_2) non-atomic then it is (\underline{B}_1 , \underline{B}_3) nonatomic as well.

<u>Proof</u>: For any measure λ , from the definition of d_{λ} it is clear that d_{λ} (\underline{B}_{1} , \underline{B}_{2}) $\leq d_{\lambda}$ (\underline{B}_{1} , \underline{B}_{3}) $\leq \frac{1}{2}$. Now, an application of theorem 2.2 gives us the result.

The next theorem helps us to generalize the result that μ is nonatomic on $\underline{\mathbb{A}}$ if and only if it is nonatomic on a countably generated sub σ -algebra $\underline{\mathbb{C}}$ of $\underline{\mathbb{A}}$. (See [1]).

Theorem 2.4 Let \underline{B} and \underline{C} be sub σ -algebras of \underline{A} . Let μ be such that (i) μ is nonatomic on \underline{C} and (ii) under μ , the sub σ -algebras \underline{B} and \underline{C} are independent. Then μ is (\underline{B} , \underline{B} \underline{V} \underline{C}) nonatomic where \underline{B} \underline{V} \underline{C} stands for the σ -algebra generated by \underline{B} and \underline{C} .

<u>Proof</u>: To prove this theorem we assume that μ is not (B, B V C) nonatomic and arrive at a contradiction.

Observe that the collection

$$\underline{\underline{D}} = \left\{ \begin{array}{cccc} (\overset{m}{\underline{D}}) & B_{j}C_{j} & \vdots & B_{j} & \epsilon & \underline{\underline{B}}, & C_{j} & \epsilon & \underline{\underline{C}}, & 1 \leq j \leq m \\ & & & & & \\ & & & & & \\ & & & & & \\ & & & & & \\ & & & & & \\ & & & & & \\ & & & & & \\ & & & & \\ & & & & \\ & & & & \\ & & & & \\ & & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & & \\ & & \\ & & & \\ & \\ & & \\ & & \\ & & \\ & & \\ & & \\ & & \\ & & \\ & & \\ & & \\ & & \\ & & \\$$

is an algebra generating $\underline{\underline{B}} \vee \underline{\underline{C}}$. Let \underline{A}_0 be a $(\underline{\underline{B}}, \underline{\underline{B}} \vee \underline{\underline{C}})(\mu)$ atom. Let $0 < \epsilon = P(\underline{A}_0)$. Choose and fix

(a) (b) (b) (c) (c)

(S₁) For any $A \in \underline{B} \ V \subseteq$, there is $B \in \underline{B}$ such that $\mu[(A \triangle B) \cap \bigcup_{j=1}^{m_0} B_j C_j] \le \epsilon/4$.

Since A is a (B, B V C)(μ) atom, given A there is B ϵ B such that μ [(A Δ B)() A] = 0. Therefore,

$$\mu[(A \triangle B) \cap \bigcup_{j=1}^{m} B_{j}C_{j}] \leq \mu[(A \triangle B) \cap A_{j}] + \mu[(A \triangle B) \cap (\bigcup_{j=1}^{m} B_{j}C_{j} - A_{j})]$$

$$\leq 0 + \mu[(\bigcup_{j=1}^{m} B_{j}C_{j}) \triangle A_{j}] \leq \varepsilon/4.$$

Hence (S1)

(S2) For $1 \le j \le m_o$, choose $C_j^o \subset C_j$ such that $\mu(C_j^o) = \frac{\mu(C_j)}{2}$. Such a choice is possible as μ is nonatomic on \underline{C} . Take $D_o = \bigcup_{j=1}^m B_j C_j^o$. By (S1) there is $B_o \in \underline{B}$ with $\mu[(D_o \triangle B_o) \subset C_j^o] \subseteq E_j^o$.

(S3)
$$\mu[(D_0^{\Delta} B_0)] = \frac{\mu((D_0^{\Delta} B_j C_j))}{2}$$
. We present

the proof below.

$$(D_{o} \triangle B_{o}) (\Box) (\underset{j=1}{\overset{m_{o}}{\bigcirc}} B_{j} C_{j}) = [((\underbrace{\overset{m_{o}}{\bigcirc}} B_{j} C_{j}^{\circ}) \triangle B_{o}) (\Box) (\underbrace{\overset{m_{o}}{\bigcirc}} B_{j} C_{j}^{\circ})]$$

$$= [(\underbrace{\overset{m_{o}}{\bigcirc}} B_{j} C_{j}^{\circ}) - B_{o}] (\underline{)} [B_{o} (\Box) (\underbrace{\overset{m_{o}}{\bigcirc}} B_{j} (C_{j} - C_{j}^{\circ}))]$$

$$(using the fact that $B_{j} C_{j} (\Box) B_{j} C_{j} = \emptyset \text{ if } j \neq j').$$$

So,
$$\mu[(D_0 \triangle B_0)(\vec{c})] = \mu[((C_j B_j C_j))] = \mu[((C_j B_j C_j)) - B_0] + \mu[B_0 C_j C_j C_j C_j)]$$

$$= \sum_{j=1}^{m} \mu(C_j^0 B_j B_0^c) + \sum_{j=1}^{m} \mu[(C_j - C_j^0) B_j B_0]$$
(using the fact that $B_j C_j (\vec{c}) B_j C_j = \emptyset$ if $j \neq j'$)
$$= \sum_{j=1}^{m} \mu(C_j^0) \mu(B_j B_0^c) + \sum_{j=1}^{m} \mu(C_j - C_j^0) \mu(B_j B_0)$$
(because of independence of B_0 and C_0)
$$= \sum_{j=1}^{m} \mu(C_j^0) \mu(B_j B_0^c) + \sum_{j=1}^{m} \mu(C_j - C_j^0) \mu(B_j B_0^c)$$

$$= \frac{1}{2} \begin{bmatrix} \sum_{j=1}^{m_o} \mu(C_j) \mu(B_j B_o^c) + \sum_{j=1}^{m_o} \mu(C_j) \mu(B_j B_o) \end{bmatrix}$$

(by the choice of G_{j}^{o} , $1 \leq j \leq m_{o}$)

$$= \frac{1}{2} \sum_{j=1}^{m_0} \mu(C_j) \mu(B_j) = \frac{1}{2} \sum_{j=1}^{m_0} \mu(C_jB_j)$$

$$= \frac{1}{2} \mu \left(\begin{array}{c} m_{o} \\ j=1 \end{array} \right) B_{j} C_{j}.$$

Hence (S3).

(S4) By (S3) and (S2)
$$\frac{\mu(\binom{n}{2}) B_j C_j}{2} \leq \epsilon/4. S_0,$$

$$\mu(\binom{n}{2}) B_j C_j \leq \epsilon/2.$$

On the other hand $\mu((\underbrace{j}_{j=1}^{m_0})B_jC_j) \triangle A_0) \leq \epsilon/4$. This implies that

 $\mu(\binom{m}{j}) \xrightarrow{B} \xrightarrow{C}_{j} \ge \varepsilon - \varepsilon/4 = \frac{3\varepsilon}{4}$. A contradiction! Hence the theorem.

Lemma 2.5 Let μ be $(\underline{\mathbb{B}}, \underline{\mathbb{A}})$ nonatomic. Let $A_0 \in \underline{\mathbb{A}}$ be such that $P_{\mu}(A_0 | \underline{\mathbb{B}}) = \varepsilon > 0$. Then there is $A_1 \in \underline{\mathbb{A}}$ $A_1 \subseteq A_0$ with $P_{\mu}(A_1 | \underline{\mathbb{B}}) = \varepsilon/2$.

Therefore, P_{μ} ($\Lambda_{1}|\underline{B}$) = $\epsilon/2$. Hence the lemma.

 $\Rightarrow \mu(A_1 \cap B) = \frac{1}{2} \mu(B)$

Theorem 2.6 μ is (B, A) nonatomic if and only if there is a countably generated sub σ -algebra $\mathcal L$ of A such that, under μ B and $\mathcal L$ are independent and μ is nonatomic on $\mathcal L$.

Proof: Let \underline{G} be a sub σ -algebra of \underline{A} such that u is nonatomic on \underline{G} and under μ , \underline{B} and \underline{G} are independent. By theorem 2.4 μ is $(\underline{B},\underline{B}\ V\ \underline{C})$ nonatomic. Since $\underline{B}\ V\ \underline{C}$ \underline{A} , by theorem 2.3 μ is $(\underline{B},\ \underline{A})$ nonatomic. This completes the proof of 'if' part.

For the proof of 'only if' part, we obtain a collection

$$A_{11}$$
 A_{12} A_{21} A_{22} A_{23} A_{24} (*)
 A_{n1} A_{n2} A_{n3} ... A_{n2}^{n}

of elements of \underline{A} as follows $A_{11} \in \underline{A}$ with $P_{\mu}(A_{11}|\underline{B}) = \frac{1}{2}$ and $A_{12} = \widehat{\bigcap} - A_{11}$; note that $P_{\mu}(A_{12}|\underline{B}) = \frac{1}{2}$. After having obtained the sets in the nth row of (*) satisfying the conditions that $\{A_{nk}, 1 \le k \le 2^n\}$ is a partition of $\widehat{\bigcap}$ by \underline{A} -sets and for any $1 \le k \le 2^n$, $P_{\mu}(A_{nk}|\underline{B}) = 1/2^n$, we obtain the sets in the (n+1)th row as follows. $A_{n+1,1} \subseteq A_{n1}$ with $P_{\mu}(A_{n+1,1}|\underline{B}) = \frac{1}{2^{n+1}}$ and $A_{n+1,2} = A_{n1} - A_{n+1,1}$. More generally for $1 \le k \le 2^n$, $A_{n+1,2k-1} \subseteq A_{nk}$ with $P_{\mu}(A_{n+1,2k-1}|\underline{B}) = \frac{1}{2^{n+1}}$ and $A_{n+1,2k} = A_{nk} - A_{n+1,2k-1} \subseteq A_{nk}$ with $P_{\mu}(A_{n+1,2k-1}|\underline{B}) = \frac{1}{2^{n+1}}$ and $A_{n+1,2k} = A_{nk} - A_{n+1,2k-1} \subseteq A_{nk}$ with $P_{\mu}(A_{n+1,2k-1}|\underline{B}) = \frac{1}{2^{n+1}}$ and $A_{n+1,2k} = A_{nk} - A_{n+1,2k-1} \subseteq A_{nk}$ with $P_{\mu}(A_{n+1,2k-1}|\underline{B}) = \frac{1}{2^{n+1}}$ is a partition of $\widehat{\bigcap}$ by \underline{A} -sets and $P_{\mu}(A_{n+1,k}|\underline{B}) = \frac{1}{2^{n+1}}$.

Define $\underline{\mathcal{C}}$ to be the σ -algebra generated by $\{\Lambda_{nk}, 1 \leq k \leq 2^n\}$ and $n \geq 1\}$. Since for $1 \leq k \leq 2^n$, $\mu(\Lambda_{nk}) = \frac{1}{2^n}$ it is clear that μ -measure of any atom of $\underline{\mathcal{C}}$ is zero; also $\underline{\mathcal{C}}$ is countably generated. Hence μ is nonatomic on $\underline{\mathcal{C}}$. Observe that the collection of sets $\{D:D$ is a union of some sets in the nth row of (*) for some $n \geq 1\}$ is an algebra generating $\underline{\mathcal{C}}$. So, to check the independence of $\underline{\mathcal{C}}$ and $\underline{\mathcal{C}}$ under μ it is enough to check that $\mu(D(\overline{\mathcal{C}})) = \mu(D) \cdot \mu(B)$ for any $B \in \underline{\mathcal{C}}$ and D from the above collection. Let $D_0 = \frac{\mathbb{C}}{k \in J} \cdot \Lambda_{n_0}^k$ where J is a subset of $1, 2, \ldots, 2^n$.

For any
$$B \in \underline{B}$$
, $\mu(D_0)$ $(\overline{D}) = \sum_{k \in J} \mu(A_{n_0k})$ $(\overline{D}) = \sum_{k \in J} \prod_{B} I_{A_{n_0k}} d\mu$

$$= \sum_{k \in J} \prod_{B} P_{\mu} (A_{n_0k}|\underline{B}) d\mu$$

$$= \sum_{k \in J} \frac{1}{2^{n_0}} \mu(B)$$

$$= \mu(B) \cdot \mu(D_0) \cdot$$

Thus \underline{B} and \underline{C} are independent under μ . This completes the proof of 'only if' part and hence the theorem.



3. The Product Problem

Let $(\underline{\zeta})_1$, \underline{A}_1 , μ_1) and $(\underline{\zeta})_2$, \underline{A}_2 , μ_2) be probability spaces. Let \underline{B}_1 ($\underline{\zeta}$) \underline{A}_1 i = 1,2. In this section we investigate the problem of $\mu = \mu_1 \times \mu_2$ being $(\underline{B}_1 (\underline{\zeta}) \underline{B}_2, \underline{A}_1 (\underline{\zeta}) \underline{A}_2)$ nonatomic.

Theorem 3.1 μ is $(\underline{\mathbb{B}}_1(\overline{x}),\underline{\mathbb{B}}_2,\underline{\mathbb{A}}(\overline{x}),\underline{\mathbb{A}}_2)$ nonatomic if and only if either μ_1 is $(\underline{\mathbb{B}}_1,\underline{\mathbb{A}}_1)$ nonatomic or μ_2 is $(\underline{\mathbb{B}}_2,\underline{\mathbb{A}}_2)$ nonatomic.

<u>Proof</u>: If there are sets A_1 and A_2 such that A_1 is a $(\underline{B}_1, \underline{A}_1)$ (μ_1) atom, i = 1, 2, then clearly $A_1 \times A_2$ is a $(\underline{B}_1(\overline{x}) \underline{B}_2, \underline{A}_1(\overline{x}) \underline{A}_2)$ $(\mu_1 \times \mu_2)$ atom. This proves the 'only if' part.

Let μ_1 be $(\underline{\mathbb{B}}_1, \underline{\mathbb{A}}_1)$ nonatoric. We shall show that $\mu_1 \times \mu_2$ is $(\underline{\mathbb{B}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{B}}_2, \underline{\mathbb{A}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{A}}_2)$ nonatomic. By Theorem 2.6 there is $\underline{\mathbb{S}}_1 (\underline{\mathbb{C}}, \underline{\mathbb{A}}_1)$ such that μ_1 is nonatomic on $\underline{\mathbb{S}}_1$ and under μ_1 the σ -algebras $\underline{\mathbb{S}}_1$ and $\underline{\mathbb{B}}_1$ are independent. It is easy to see that $\mu_1 \times \mu_2$ is nonatomic on $\underline{\mathbb{S}}_1 (\overline{\mathbf{x}}) \{(\emptyset, (\widehat{\mathbb{C}}_2))\}$ and under $\mu_1 \times \mu_2$, $\underline{\mathbb{S}}_1 (\overline{\mathbf{x}}) \{\emptyset, (\widehat{\mathbb{C}}_2)\}$ and $\underline{\mathbb{B}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{B}}_2$ are independent. Hence by Theorem 2.4 the measure $\mu_1 \times \mu_2$ is $(\underline{\mathbb{B}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{B}}_2, \underline{\mathbb{B}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{B}}_2 \vee \underline{\mathbb{S}}_1 (\overline{\mathbf{x}}) \{\emptyset, (\widehat{\mathbb{C}}_2)\}$ nonatomic. Now, an application of Theorem 2.3 gives us that $\mu_1 \times \mu_2$ is $(\underline{\mathbb{B}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{B}}_2, \underline{\mathbb{A}}_1 (\overline{\mathbf{x}}) \underline{\mathbb{A}}_2)$ nonatomic. This completes the proof of 'if' part and hence the theorem.

Remark: The above theorem generalizes the theorem of [1] for product measures.

Let λ be a measure on $(\underbrace{\zeta}_1 \times \underbrace{\zeta}_2, \underbrace{A}_1(\underline{x}) \underline{A}_2)$. Let λ_i be the marginal on $(\underbrace{\zeta}_i, \underline{A}_i)$, i = 1, 2. Then it is known that if one of the λ_i 's is nonatomic then so is λ (See [1]). This however does not generalize as shown by the example given below.

Example: For i = 1, 2 let $\bigcap_{i} = [0, 1]$ and $\underline{A}_{i} = \text{the Borel}$ σ -algebra on [0,1]. Take $\underline{B}_{1} = \{\emptyset, \bigcap_{1}\}$ and $\underline{B}_{2} = \underline{A}_{2}$. Let λ be the uniform measure on the diagonal of $\bigcap_{1} \times \bigcap_{2}$. Then it is clear that $\lambda_{1} = \lambda_{2} = \text{the Lebesgue measure on } [0,1]$. So λ_{1} is $(\underline{B}_{1}, \underline{A}_{1})$ nonatomic. However λ is not $(\underline{B}_{1}(\overline{x}) \underline{B}_{2}, \underline{A}_{1}(\overline{x}) \underline{A}_{2})$ nonatomic since the diagonal of $(\underline{A}_{1} \times \widehat{A}_{2})$ is a $(\underline{B}_{1}(\overline{x}) \underline{B}_{2}, \underline{A}_{1}(\overline{x}) \underline{A}_{2})(\lambda)$ atom.

4. A generalization of Liapounoff's theorem

In this section we obtain a generalization of the classical Liapounoff's theorem. Our proof is essentially that of Lindenstrauss' of the classical version [13] with suitable modifications.

Theorem 4.1 (Liapounoff's Theorem) Let $\mu_1, \mu_2, \ldots, \mu_n$ be (B, A) nonatomic measures. Equip I_{∞} (C), B, μ_1) with the

w* topology (that is, the weak topology induced by $L_1(\underline{C}, \underline{B}, \mu_1)$); let $L_{\infty}(\mu_1)$ stand for this space. The set $(P_{\mu_1}(A|\underline{B}), P_{\mu_2}(A|\underline{B}), \dots, P_{\mu_n}(A|\underline{B}))/A \in \underline{A}$ is a closed convex subset of $L_{\infty}(\mu_1) \times L_{\infty}(\mu_2) \times \dots \times L_{\infty}(\mu_n)$.

Proof: Let n=1. By theorem 1.1 it is clear that $\{P_{\mu_1}(A|B)/A \in A\} = \{f: 0 \le f \le 1 \text{ and } f \in I_{\infty}(\mu_1)\}$ and hence the result. We complete the proof by induction.

Let $\mu = \mu_1 + \mu_2 + \cdots + \mu_n$. Let $W = \{g: 0 \le g \le 1, g \in I_{\infty} (\Omega), \underline{A}, \mu\}$. Equip $I_{\infty} (\Omega), \underline{A}, \mu$ with the w^* -topology. Then it is known that W is a compact set. Clearly, W is convex as well. Define a mapping T from W to $I_{\infty} (\mu_1) \times I_{\infty} (\mu_2) \times \cdots \times I_{\infty} (\mu_n)$ by setting, for g in W, $T(g) = (P_{\mu_1}(g|\underline{B}), P_{\mu_2}(g|\underline{B}), \dots, P_{\mu_n}(g|\underline{B}))$. Presently we show that T is continuous. Let $\{g_{\infty}\}$ be a net in W converging to g_0 in W. Consider $h \in L_1(\Omega)$, \underline{B} , μ_1 for some I. Observing that I is in $I_1(\Omega)$, I I I I I I we have

$$\int P_{\mu_{\mathbf{i}}}(g_{\alpha}|\underline{B})h \, d\mu_{\mathbf{i}} = \int g_{\alpha} \cdot h \, d\mu_{\mathbf{i}}$$

$$= \int g_{\alpha} \cdot h \cdot \frac{d\mu_{\mathbf{i}}}{d\mu} d\mu \longrightarrow \int g_{0} h \, \frac{d\mu_{\mathbf{i}}}{d\mu} d\mu$$

$$= \int g_{0} h \, d\mu_{\mathbf{i}} = \int P_{\mu_{\mathbf{i}}}(g_{0}|\underline{B}) h \, d\mu_{\mathbf{i}}.$$

Since h and i are arbitrary it follows that T is continuous. So, T(W) is a compact (and hence closed) and convex (since T is affine) subset of $I_{\infty}(\mu_1) \times I_{\infty} (\mu_2) \times \cdots \times I_{\infty} (\mu_n)$. The proof would be complete if we could show $T(W) = \{T(I_D) \mid D \in A\}$.

Let (h_1, h_2, \ldots, h_n) ε T(W). We have to exhibit $D \varepsilon \underline{A}$ with $T(I_D) = (h_1, h_2, \ldots, h_n)$. Let $W_0 = T^{-1}(h_1, h_2, \ldots, h_n)$. W_0 is a compact convex subset of W_0 . By Krein-Milman theorem W_0 has extreme points. Let g be an extreme point of W_0 . We shall prove that $g = I_D$ a.s. (μ) for some $D \varepsilon \underline{A}$.

Suppose not. Then for some $\varepsilon > 0$, $\mu(\{\varepsilon \le g \le 1 - \varepsilon\}) > 0$. This implies that for some i, say i = 1, $\mu_i(\{\varepsilon \le g \le 1 - \varepsilon\}) > 0$; that is $\mu_i(\{\varepsilon \le g \le 1 - \varepsilon\}) > 0$. Denote $\{\varepsilon \le g \le 1 - \varepsilon\}$ by Z. Choose an Δ measurable subset A of Z such that for every $B \in B$, $\mu_i[A \land (B \land Z)] > 0$ and $\mu_i[A \land (B \land Z)] > 0$. (Such a choice is possible since Z is not a $(B, \Delta)(\mu_i)$ atom). By induction hypothesis, we can find Δ measurable sets B and C satisfying (i) $B \land A$, $C \land Z - A$ and (ii) for $2 \le i \le n$, $2P_{\mu_i}(B|B) = P_{\mu_i}(A|B)$ and $2P_{\mu_i}(C|B) = P_{\mu_i}(Z - A|B)$.

Choose and fix versions of $P_{\mu_1}(B|B)$, $P_{\mu_1}(A|B)$, $P_{\mu_1}(C|B)$ and $P_{\mu_1}(Z-A|B)$. Denote by x and y the functions $2P_{\mu_1}(B|B) - P_{\mu_1}(A|B) \text{ and } 2P_{\mu_1}(C|B) - P_{\mu_1}(Z-A|B) \text{ respectively.}$

Define two B-measurable functions s and t as follows.

 $s = \varepsilon$ and $t = \varepsilon$ on $\{x = 0, y = 0\}$,

 $s = \varepsilon$ and t = 0 on $\{x = 0, y \neq 0\}$,

s = 0 and $t = \varepsilon$ on $\{x \neq 0, y = 0\}$,

 $s = t \cdot \frac{y}{x}$ and $t = \epsilon$ on $\{0 < |\frac{y}{x}| \le 1\}$ and

 $s = \varepsilon$ and $t = s \cdot \frac{x}{y}$ on $\{|\frac{y}{x}| > 1\}$.

Clearly $|s(w)| \le \varepsilon$ and $|t(w)| \le \varepsilon$ for all $w \in \Gamma$ and sx = ty. Define a function $h by h = s(2I_B - I_A) + t(I_{Z-A} - 2I_C)$.

1°. For $1 \le i \le n$, $P_{\mu_i}(h|\underline{B}) = 0$.

If $2 \le i \le n$, this is clear from the choice of sets B and C; for i = 1, this is by the choice of functions s and t.

 2° . g + h and g - h are in W_{\bullet}

This is because $|h| \le g \le 1 - |h|$ on \bigcirc ; so $0 \le g + h \le 1$ and $0 \le g - h \le 1$.

- 3° . g + h and g h are in W_{\circ} . This is immediate from 1° and 2° .
- 4° . $\mu_{1}(h \neq 0) > 0$.

Suppose not. Then h = 0 on Z a.s. (μ_1) . Note that

h = s or -s on A and h = t or -t on Z - A. So, s = 0 on A a.s. (μ_1) and t = 0 on Z - A a.s. (μ_1) . Then, from the definition of s and t we have

A \subset $\{s = 0\} = \{x \neq 0, y = 0\}$ a.s. (μ_1) and $z - A \subset \{t = 0\} = \{x = 0, y \neq 0\}$ a.s. (μ_1) .

Now, $\mu_{1}(A \triangle [\{s = 0\}](B)) = \mu_{1}[(A \triangle \{s = 0\})(B)] = \mu_{1}(\{s = 0\}](B) = 0$

A contradiction to the choice of A since $\{s=0\}$ is a B measurable set. Thus 4° .

Now, from 4° it is clear that $\mu(h \neq 0) > 0$. In view of 3° , therefore, it follows that g is not an extreme point of W. This contradiction concludes the proof of the theorem.

Remark: In the above we have shown more than what we set out to prove. We have shown that the set $\{ (P_{\mu_1}(A|B), P_{\mu_2}(A|B), \dots, P_{\mu_n}(A|B)) / A \in \underline{A} \} \text{ is convex even when the 'weights' are \underline{B}-measurable functions. The paragraph below is intended to amplify this.}$

Let $\mu_1, \, \mu_2, \dots, \mu_n$ be ($\underline{\mathbb{B}}$, $\underline{\mathbb{A}}$) nonatomic measures. Let \mathbb{W} and \mathbb{T} be as in theorem 4.1. Let $A_1, \, A_2$ be elements of $\underline{\mathbb{A}}$. Let $g_1, \, g_2$ be two $\underline{\mathbb{B}}$ -measurable functions with $0 \leq g_i \leq 1$, i = 1, 2 and $g_1 + g_2 \leq 1$. Clearly $g_1 = 1$, $a_1 + g_2 = 1$, $a_2 \in \mathbb{W}$. Now

$$P_{\mu_{\hat{\mathbf{j}}}}(g_{1} I_{A_{1}} + g_{2} I_{A_{2}} | \underline{\underline{B}}) = g_{1} P_{\mu_{\hat{\mathbf{j}}}}(I_{A_{1}} | \underline{\underline{B}}) + g_{2} P_{\mu_{\hat{\mathbf{j}}}}(I_{A_{2}} | \underline{\underline{B}}),$$

$$(as \ g_{1} \text{ and } g_{2} \text{ are } \underline{\underline{B}} \text{ measurable}),$$

$$j = 1, 2, \dots n.$$

So,
$$T(g_{\underline{1}} I_{\underline{A}_{\underline{1}}} + g_{\underline{2}} I_{\underline{A}_{\underline{2}}}) = g_{\underline{1}} T(I_{\underline{A}_{\underline{1}}}) + g_{\underline{2}} T(I_{\underline{A}_{\underline{2}}})$$
. Since
$$T(W) = \{T(I_{\underline{D}})/D \in \underline{\underline{A}}\} \text{ we have that}$$
$$g_{\underline{1}} T(I_{\underline{A}_{\underline{1}}}) + g_{\underline{2}} T(I_{\underline{A}_{\underline{2}}}) \in \{T(I_{\underline{D}}) / D \in \underline{\underline{D}}\}.$$

As a consequence of the above theorem we have the following generalization of theorem 2.6.

Theorem 4.2 Let μ_1 , μ_2 , ..., μ_n be ($\underline{\underline{B}}$, $\underline{\underline{A}}$) nonatomic measures. Then there is a countably renerated sub σ -algebra $\underline{\underline{C}}$ of $\underline{\underline{A}}$ such that each μ_1 is nonatomic on $\underline{\underline{C}}$ and under each μ_1 , $\underline{\underline{B}}$ and $\underline{\underline{C}}$ and independent.

<u>Proof</u>: Use theorem 4.1 and imitate the construction of $\underline{\mathcal{L}}$ in theorem 2.6. We omit the details.

5. A generalization of a result of Halmos

Theorem 5.1. Let μ_1 and μ_2 be two nonatomic measures on ((), A). Let for some r with 0 < r < 1,

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 $\{A \in \underline{A} / \mu_1(A) = r \} = \{A \in \underline{A} / \mu_2(A) = r \} \cdot \text{ Then } \mu_1 = \mu_2$

This result is implicit in Halmos's proof of Liapounoff's theorem [8], as has been pointed out by Bolkar. For a proof of this result see corallary 1.23 of [6] or theorem 1.5.1 of [3]. Below, we obtain a generalization of this result.

Theorem 5.2 Let μ_1 and μ_2 be two ($\underline{\underline{B}}$, $\underline{\underline{A}}$) nonatomic measures. Let f_0 be a $\underline{\underline{B}}$ measurable function with the property $\mu_1(\{0 < f_0 < 1\}) = 1$. Further, let

$$\{A \in \underline{A} / P_{\mu_{1}}(A|\underline{B}) = f_{o}\} = \{A \in \underline{A} / P_{\mu_{2}}(A|\underline{B}) = f_{o}\}.$$

Then $\mu_1 = \mu_2$

The proof is completed with the help of two lemmas.

Lemma 5.3 Assume the hypothesis of theorem 5.2. Then $\mu_1 = \mu_2$ on Δ . Also for any $B_0 \in \underline{B}$.

$$\{A \in \underline{A} / P_{\mu_1}(A|\underline{B}) = f_o \cdot I_{B_o}\} = \{A \in \underline{A} / P_{\mu_2}(A|\underline{B}) = f_o \cdot I_{B_o}\} \cdot$$

Proof: Let $A \in \underline{A}$ with $\mu_1(A) = 0$. Choose $C \in \underline{A}$ with $P_{\mu_1}(C|\underline{B}) = f_0$. (Such a choice is possible in view of theorem 1.1). Then $P_{\mu_1}(C(\underline{A})) = f_0$ and $P_{\mu_1}(C(\underline{A})) = f_0$. So, $P_{\mu_2}(C(\underline{A})) = f_0$ and $P_{\mu_2}(C(\underline{A})) = f_0$. Hence, $P_{\mu_2}(C(\underline{A})) = f_0$.

i.e. $P_{\mu_2}(A|B) = 0$. This implies $\mu_2(A) = 0$. Thus $\mu_2 << \mu_1$. Interchanging the roles of μ_1 and μ_2 in the above we conclude $\mu_1 << \mu_2$. Thus $\mu_1 \equiv \mu_2$ on \underline{A} .

Fix $B_0 \in \underline{B}$. Let $A_0 \in \underline{A}$ be such that $P_{\mu_1}(A_0 | \underline{B}) = f_0 I_{B_0}$.

To show $P_{\mu_2}(A_0 | \underline{B}) = f_0 I_{B_0}$

$$\mu_{1}(\Lambda_{0} \stackrel{\frown}{}) \stackrel{B^{c}}{=} = \int_{B^{c}_{0}} I_{A_{0}} d\mu_{1} = \int_{B^{c}_{0}} P_{\mu_{1}}(\Lambda_{0} | \underline{B}) d\mu_{1}$$

$$= \int_{B^{c}_{0}} f_{0} \cdot I_{B_{0}} d\mu_{1} = 0.$$

So, μ_2 (Λ_0 ($\tilde{}$) B_0^c) = 0 (as $\mu_1 = \mu_2$). Choose $C \in \underline{A}$ with $P_{\mu_1}(C|\underline{B}) = f_0$ and $P_{\mu_2}(C|\underline{B}) = f_0$. (Such a choice is possible in view of the remark following theorem 4.1). Let D = C ($\tilde{}$) B_0^c . Then $P_{\mu_1}(D|\underline{B}) = f_0 \cdot I_{B^c}$ and $P_{\mu_2}(D|\underline{B}) = f_0 \cdot I_{B^c}$. Now,

$$P_{\mu_{1}}(A_{o}(\underline{\ })D|\underline{B}) = P_{\mu_{1}}(A_{o}|\underline{B}) + P_{\mu_{1}}(D|\underline{B})$$

$$(\bullet, \bullet, \mu_{1}(A_{o}(\underline{\ })D) \leq \mu_{1}(A_{o}(\underline{\ })B_{o}) = 0)$$

$$= f_0 I_{B_0} + f_0 I_{B_0} = f_0.$$

 $\cdot \cdot \cdot P_{\mu_{2}} (A_{0} (\underline{\ }) D|\underline{B}) = f_{0} \cdot$

i.e. $P_{\mu_2}(A_0|\underline{B}) + P_{\mu_2}(D|\underline{B}) = f_0 (\cdot \cdot \mu_2(A_0|\underline{C}) D) \le \mu_2(A_0|\underline{C}) B_0^C = 0$

$$i_{\bullet}e_{\bullet} \quad P_{\mu_{2}}(A_{0}|\underline{B}) + f_{0} I_{B_{0}} = f_{0}.$$

Hence $P_{\mu_2}(A_0|\underline{B}) = f_0 I_{B_0}$. Interchanging the roles of μ_1 and μ_2 the proof of the lemma can be completed.

Lemma 5.4 Assume the hypothesis of theorem 5.2. Let g be a B measurable function with $0 \le g \le 1$. Then

(i)
$$\{A \in \underline{A} / P_{\mu_1}(A | \underline{B}) = g I_{\{g = f_o\}}\} = \{A \in \underline{A} / P_{\mu_2}(A | \underline{B}) = g I_{\{g = f_o\}}\},$$

(ii)
$$\{A \in \underline{A} / P_{\mu_{\underline{1}}}(A|\underline{B}) = g I_{\{g > f_{\underline{0}}\}} = \{A \in \underline{A} / P_{\mu_{\underline{2}}}(A|\underline{B}) = g I_{\{g > f_{\underline{0}}\}}\}$$
 and

(iii)
$$\{ \Lambda \in \underline{\Lambda} / P_{\mu_1}(\Lambda | \underline{B}) = g I_{\{g < f_o\}} \} =$$

$$\{ \Lambda \in \underline{\underline{\Lambda}} / P_{\mu_2}(\Lambda | \underline{B}) = g I_{\{g < f_o\}} \}.$$

Proof. Since
$$g I_{g = f_o} = f_o I_{g = f_o}$$
 and $g = f_o$

is a B measurable set an application of lemma 5.3 proves (i).

Let $A_0 \in \underline{A}$ be such that $P_{\mu_1}(A_0|\underline{B}) = g I_{\{g > f_0\}}$. To show $P_{\mu_2}(A_0|\underline{B}) = g I_{\{g > f_0\}}$. Since $\mu_1(\{f_0 > 0\}) = 1$, the function $\frac{f_0}{g}$ can be meaningfully talked of. Take $h = \frac{f_0}{g}$ Al. By the remark following theorem 4.1, there is a D ϵ Δ with

$$(P_{\mu_{1}}(D|\underline{B}), P_{\mu_{2}}(D|\underline{B})) = (h P_{\mu_{1}}(A_{o}|\underline{B}), h P_{\mu_{2}}(A_{o}|\underline{B})).$$

$$Now, P_{\mu_{1}}(D|\underline{B}) = h P_{\mu_{1}}(A_{o}|\underline{B})$$

$$= h g I_{g} > f_{o};$$

$$= \frac{f_{o}}{g} g I_{g} > f_{o};$$

$$= f_{o} I_{g} > f_{o};$$

$$= f_{o} I_{g} > f_{o};$$

$$(as h = \frac{f_{o}}{g} \text{ on the set})$$

$$= f_{o} I_{g} > f_{o};$$

. By lemma 5.3,
$$P_{\mu_2}(D|\underline{B}) = f_0 I_{\{g > f_0\}}$$

i.e. $h P_{\mu_2}(A_0|\underline{B}) = f_0 I_{\{g > f_0\}}$ (*)

Observe that since $\mu_2([f_0 > 0]) = 1$, we have $\mu_2([h = 0]) = 0$. Now, r.h.s. of (*), and consequently h $P_{\mu_2}(A_0|B)$, is 0 on $[g \le f_0]$. So, $P_{\mu_2}(A_0|B) = 0$ a.s. (μ_2) on $\{g \le f_0\}$. Again, r.h.s. of (*), and consequently h $P_{\mu_2}(A_0|B)$, is f_0 on $\{g > f_0\}$

i.e.
$$\frac{f_0}{g} P_{\mu_2}(A_0|B) = f_0 \text{ on } \{g > f_0\} \text{ (as } h = \frac{f_0}{g} \text{ on } \{g > f_0\} \text{)}$$

i.e.
$$P_{\mu_2}(A_0|B) = g$$
 a.s. (μ_2) on $\{g > f_0\}$, since $\mu_2(\{f_0 > 0\}) = 1$.

Thus
$$P_{\mu_2}(A_0|\underline{B}) = g I_{g} > f_0$$

The other part of (ii) is similarly proved . Hence (ii).

Let $A_0 \in \underline{A}$ be such that $P_{\mu_1}(A_0 | \underline{B}) = g I_{g} < f_0$.

Choose Dε A with

$$(P_{\mu_{1}}(D|\underline{B}), P_{\mu_{2}}(D|\underline{B})) = (\frac{1-f_{0}}{1-g} \wedge 1) (P_{\mu_{1}}(A_{0}|\underline{B}), P_{\mu_{2}}(A_{0}|\underline{B})) + I_{\{g < f_{0}\}} \frac{f_{0}-g}{1-g} (1, 1).$$

Verifying that there is such a D is routine (using the remark following theorem 4.1) and is omitted.

Now,
$$P_{\mu_{1}}(D|B) = (\frac{1-f_{0}}{1-g} \wedge 1) P_{\mu_{1}}(A_{0}|B) + I_{\{g < f_{0}\}} \cdot \frac{f_{0}-g}{1-g}$$

$$= (\frac{1-f_{0}}{1-g} \wedge 1) g I_{\{g < f_{0}\}} + I_{\{g < f_{0}\}} \cdot \frac{f_{0}-g}{1-g}$$

$$= f_{0} I_{\{g < f_{0}\}} \cdot \frac{f_{0}-g}{1-g}$$

So, by lemma 5.3,
$$P_{\mu_2}(D|\underline{B}) = f_0 I_{g < f_0}$$

i.e.
$$(\frac{1-f_0}{1-g} \land 1) P_{\mu_2}(A \mid B) + I_{\{g < f_0\}} \frac{f_0-g}{1-g} = f_0 I_{\{g < f_0\}}$$
 (**)

$$\frac{1-f_0}{1-g} P_{\mu_2}(A_0 | \underline{B}) + \frac{f_0-g}{1-g} = f_0$$

or
$$(1-f_0) P_{\mu_2}(A_0, B) + f_0-g = f_0(1-g)$$
 (note that since $f_0 < 1$, $g < f_0 \implies g < 1$)

or
$$(1-f_0) P_{\mu_2}(A_0|B) = g(1-f_0)$$

or
$${}^{\mathbf{P}}\mu_{\mathbf{C}}(\mathbf{A}_{\mathbf{O}}|\mathbf{B}) = \mathbf{g}$$

i.e.
$$P_{\mu_2}(A_0|B) = g$$
 a.s. (μ_2) on $\{g < f_0\}$.

On
$$fg \ge f_0$$
 (**) reduces to

$$(\frac{1-f_0}{1-g} \quad \Lambda \quad 1) \quad P_{\mu_2}(\Lambda_0 \mid \underline{B}) = 0$$

or
$$P_{\mu_2}(\Lambda_0 | \underline{B}) = 0$$

i.e.
$$P_{\mu_2}(\Lambda_0|\underline{B}) = 0$$
 a.s. (μ_2) on $\{g \ge f_0\}$.

Thus
$$P_{\mu_2}(A_0|B) = g I_{g} < f_0$$
.

Interchanging the roles of μ_1 and μ_2 the proof of (iii) is completed.

Hence the lemma.

Proof of theorem 5.2: Let
$$A_0 \in \underline{A}$$
. Let $g = P_{\mu_1}$ ($A_0 \mid \underline{B}$). Then

$$A_0 = (A_0 () \{g < f_0\}) () (A_0 () \{g = f_0\}) () (A_0 () \{g > f_0\}).$$

Now,

So, by lemma 5.4,

$$\begin{array}{lll} P_{\mu_{2}}(A_{0} & \frown) & \{g < f_{0}\} \mid B\} = & g \cdot I_{\{g < f_{0}\}} \cdot \\ P_{\mu_{2}}(A_{0} & \frown) & \{g = f_{0}\} \mid B\} = & g \cdot I_{\{g = f_{0}\}} \quad \text{and} \\ P_{\mu_{2}}(A_{0} & \frown) & \{g > f_{0}\} \mid B\} = & g \cdot I_{\{g > f_{0}\}} \cdot \\ \end{array}$$

Hence P_{μ_2} (A) \underline{B}) = 3.

$$P_{\mu_1}(A_0|\underline{B}) = P_{\mu_2}(A_0|\underline{B}).$$
 This implies $\mu_1(A_0) = \mu_2(A_0).$

Hence the theorem.

In the above theorem the condition that $\mu_1(0 < f_0 < 1)=1$ (which is analogous to the condition that 0 < r < 1 in theorem 5.1) cannot be relaxed. For this, consider the following example.

Let () = [0, 1], \underline{A} its Borel σ -algebra. Let μ be the Lebesgue measure on $((), \underline{A})$. Let α_1 and α_2 be numbers such that $\alpha_1 \neq \alpha_2$, $0 < \alpha_i < 1$, i = 1, 2. Define measures μ_1

and
$$\mu_2$$
 by setting for $\Lambda \in \underline{\Lambda}$,

$$\mu_{1}(\Lambda) = \frac{\alpha_{1}}{2} \frac{\mu(\Lambda \cap [0, \frac{1}{4}])}{\mu([0, \frac{1}{4}])} + \frac{(1 - \alpha_{1})}{2} \frac{\mu(\Lambda \cap [\frac{1}{4}, \frac{1}{2}])}{\mu([\frac{1}{4}, \frac{1}{2}])} + \frac{\mu(\Lambda \cap [\frac{1}{4}, \frac{1}{2}])}{\mu([\frac{1}{4}, \frac{1}{2}])} + \frac{\mu(\Lambda \cap [0, \frac{1}{4}])}{\mu([0, \frac{1}{4}])} + \frac{(1 - \alpha_{2})}{2} \frac{\mu(\Lambda \cap [\frac{1}{4}, \frac{1}{2}])}{\mu([\frac{1}{4}, \frac{1}{2}])} + \frac{\mu(\Lambda \cap [\frac{1}{4}, \frac{1}{4}])}{\mu([\frac{1}{4}, \frac{1}{2}])} + \frac{\mu(\Lambda \cap [\frac{1}{4}, \frac{1}{4}])}{\mu([\frac{1}{4}, \frac{1}{4}])} + \frac{\mu(\Lambda \cap [\frac{1}{4}, \frac{1}{4}])}{\mu([\frac{1}{4}, \frac{1}{4}])}$$

Let
$$\underline{B} = \{ \emptyset, [0, \frac{1}{2}), [\frac{1}{2}, 1], () \}$$
. Clearly μ_1 and μ_2 are $(\underline{B}, \underline{A})$ nonatomic measures. Now, for any A with $\mu(A() [0, \frac{1}{2}]) = 0$,

$$P_{\mu_1}(A \mid \underline{B}) = 0$$
 on $[0, \frac{1}{2})$
= $2\mu(A(\overline{)})[\frac{1}{2}, 1]$) on $[\frac{1}{2}, 1]$, and

$$P_{\mu_2}(A \mid \underline{B}) = 0$$
 on $[0, \frac{1}{2})$
= $2 \mu (A \stackrel{\frown}{}) [\frac{1}{2}, 1])$ on $[\frac{1}{2}, 1]$; that is

$$p_{\mu_1}(A|\underline{B}) = P_{\mu_2}(A|\underline{B}).$$

Let
$$f_0 = I$$
 [$\frac{1}{2}$, 1]

$$P_{\mu_{1}}(A|\underline{B}) = f_{0} \Rightarrow \mu(A(\overline{})[0, \frac{1}{2})) = 0$$

$$\Rightarrow P_{\mu_{2}}(A|\underline{B}) = P_{\mu_{1}}(A|\underline{B})$$

$$\Rightarrow P_{\mu_{2}}(A|\underline{B}) = f_{0}.$$

Likewise
$$P_{\mu_2}(A|B) = f_0 \Rightarrow P_{\mu_1}(A|B) = f_0$$
.
However $\mu_1 \neq \mu_2$ as $\mu_1([0, \frac{1}{4}]) = \frac{\alpha_1}{2}$ and $\mu_2([0, \frac{1}{4}]) = \frac{\alpha_2}{2}$.

6. Nonatomicity in Polish Spaces

Let \subseteq be a Polish space and \triangle its Borel σ -algebra. In this section we consider the problem of characterizing those sub σ -algebras \triangle of \triangle which admit a (\triangle, \triangle) non-atomic measure; also, we give an equivalent characterization of a measure λ being (\triangle, \triangle) nonatomic in terms of the atoms of \triangle . For our study, we need the following two results.

Theorem 6.1 (Blackwell). Let $A \in \underline{A}$. Let \underline{B} be a countably generated sub σ -algebra of \underline{A} such that $A(\overline{\ })$ \underline{B} has singleton atoms. Then $A(\overline{\ })$ $\underline{B} = A(\overline{\ })$ A.

This follows as an easy consequence of Corollary 1 to theorem 3 of [4].

Theorem 6.2 (Lusin). Let \underline{B} be a countably generated sub σ -algebra of \underline{A} with every atom countable. Then we can get a sequence $\{A_n\}_{n\geq 1}$ of Borel sets such that

(1)
$$\Lambda_{i}$$
 (7) $\Lambda_{j} = \emptyset$, $i \neq j$ and $(\sum_{n=1}^{\infty} A_{n} = I)$ and

(2) each A_n contains at most one point from every atom of B.

A proof of this is available in page 335 of [10].

The next theorem gives a complete solution to the problem mentioned in the beginning of this section when the sub o-algebra considered is countably generated.

Theorem 6.3 Let \underline{B} be a countably generated sub σ -algebra of \underline{A} . Then there exists a $(\underline{B}, \underline{A})$ nonatomic measure if and only if \underline{B} has at least one uncountable atom.

Proof. Let \underline{B} have an uncountable atom, say \underline{B}_0 . By Borel isomorphism theorem \underline{B}_0 supports a nonatomic measure, say λ . Let λ_1 be the extension of λ to \underline{A} defined by $\lambda_1(A) = \lambda(A(\overline{C})B_0)$, $A \in \underline{A}$. It is clear that λ_1 is a nonatomic measure. We shall show that λ_1 is $(\underline{B}, \underline{A})$ nonatomic. Let $A \in \underline{A}$ with $\lambda_1(A) > 0$; so $\lambda_1(A(\overline{C})B_0) > 0$. Choose $A_1(A(\overline{C})B_0)$ with $\lambda_1(A_1) = \frac{1}{2}\lambda_1(A(\overline{C})B_0)$.

(A(\overline{C})B_0(\ove

 $= A (\overline{)} \{\emptyset, (\overline{)}\} \quad \text{(since } B_0 \text{ is a } \underline{B}\text{-atom})$ $= \{\emptyset, \Lambda\}.$

On the other hand $A_1 \in A \cap B_0 \cap A$ and $A_1 \cap A_1 \cap A \cap A \cap A_1 \cap A_1 \cap A \cap A_1 \cap A_1$

Let now every atom of $\underline{\beta}$ be countable. By theorem 6.2 there are Borel sets $\Lambda_1, \Lambda_2, \ldots, \Lambda_n, \ldots$ such that

(1)
$$A_j = \emptyset$$
 if $i \neq j$ and $(-)_{n=1}^{\infty} A_n = (-)_n$ and

(2) each A_n contains at most one point from every atom of \underline{B} . We now claim that A_n () $\underline{B} = A_n$ () \underline{A} , for every n. (Note that there is no measure involved). For, A_n () \underline{B} has singleton atoms (in view of property (2) mentioned above); hence by theorem 6.1, A_n () $\underline{B} = A_n$ () \underline{A} . Let λ be any measure on \underline{A} . Let n be such that $\lambda(A_n) > 0$. Then clearly A_n is a $(\underline{B}, \underline{A})$ (λ) atom. Thus no measure on \underline{A} is a $(\underline{B}, \underline{A})$ nonatomic measure. Hence the 'only if' part and the theorem.

Regarding the case of not-countably generated B we have the following two examples.

Example 1: Let \bigcirc be uncountable. Let \bigcirc B = \bigcirc B /B is Borel, either B or B^c is a countable set \bigcirc

Clearly \underline{B} is not countably generated. Let λ be a nonatomic measure on ($\underline{\Gamma}$), \underline{A}). (There exist many such). Since $\{\emptyset, \underline{\Gamma}\} = \underline{B} \pmod{\underline{N}}$ it is clear that λ is $(\underline{B}, \underline{A})$ nonatomic. Thus there are not-countably generated sub σ -algebras \underline{B} admitting $(\underline{B}, \underline{A})$ nonatomic measures.

Example 2: Let () = [0, 1]. Let A be a non-Borel universal null set. (The existence of such sets assumes the continuum hypothesis and the axim of choice. For details see pages 525 and 532 of [11]).

Define $\underline{B} = \{B \mid B \text{ is Borel, and either } B \subset A^C \text{ or } B^C \subset A^C \}$.

That \underline{B} is a σ -algebra is evident. Since \underline{A} is a non-Borel set, it follows that \underline{B} is not countably generated. Let λ be any nonatomic measure on \underline{A} . By the choice of A, we have A is λ -measurable and $\lambda(A^C) = 1$. So, $\underline{B} = \underline{\Lambda}$ (modulo \underline{N}_{λ}). Hence λ is not $(\underline{B}, \underline{A})$ nonatomic. Thus \underline{B} does not admit any $(\underline{B}, \underline{A})$ nonatomic measure. So, there are not-countably generated sub σ -algebras \underline{B} not admitting any $(\underline{B}, \underline{\Lambda})$ non-atomic measure.

Before stating our next result we need a definition.

Definition: Let \underline{B} be a countably generated sub σ -algebra of \underline{A} . As \underline{A} is said to be one-sheeted with respect to \underline{B} if it contains at most one point from each atom of \underline{B} .

Theorem 6.4 Let λ be a measure on ($(\ \ \)$, $\underline{\wedge}$). Let \underline{B} be a countably generated sub σ -algebra of $\underline{\wedge}$. Then λ is (\underline{B} , $\underline{\wedge}$) nonatomic if and only if there is no $\underline{\wedge}$ a which is one-sheeted with respect to \underline{B} and has positive λ -measure.

Proof: Let D be a (\underline{B} , \underline{A}) (λ) atom. Let $\{A_n\}_{n\geq 1}$ be a sequence of sets generating \underline{A} . Since D is a (\underline{B} , \underline{A}) (λ) atom, for each n we can get a \underline{B} -set \underline{B}_n with $\lambda((B_n \land A_n) \cap D) = 0$. Let $\Lambda = D - \binom{n}{2} \cdot \binom{n}{2}$

(A () B

... A () \underline{A} = A () \underline{B} . Since A () \underline{A} has singleton atoms, A () \underline{B} has singleton atoms. Hence A contains at most one point from each atom of \underline{B} . This completes the proof of 'if' part.

Let now A be one-sheeted with respect to \underline{B} and have positive λ -measure. Then Λ () \underline{B} is a sub σ -algebra of Λ () $\underline{\Delta}$ and has singleton atoms. Therefore, by theorem 6.1,

 Λ () $B = \Lambda$ () A; in particular, Λ is a (B, A) (Λ) atom. This completes the proof of 'only if' part and hence the theorem.

Remarks :

- (1) The concept of 'one-sheeted' subsets has been introduced by Rohlin in [22], for Lebesgue spaces.
- (2) Theorem 6.4 is not the best possible. A careful perusal of the proof of the theorem reveals that the 'if' part remains valid so long \underline{A} is a countably generated σ -algebra with singleton atoms; no assumption of topological nature need be made on \underline{C} . However the same thing cannot be said of the 'only if' part. Below we present an example to elucidate this.

So, $\lambda_1((A_1 \triangle B) (\)(A (\)D))>0$ for all $B \in \underline{B}$.

i.e. $\lambda_1((A_1 D \triangle B D) (\)AD)>0$ for all $B \in \underline{B}$.

Hence the traces of $D (\)\underline{A}$ and $D (\)\underline{B}$ do not coincide (upto \underline{N}_1 sets) on $D (\)A$. So, $D (\)A$ is not a ($D (\)\underline{B}$, $D (\)\underline{A}$) (λ_1) atom. Thus λ_1 is ($D (\)\underline{B}$, $D (\)\underline{A}$) nonatomic.

Now let us specialize. Take () = [0, 1] \times [0, 1], \underline{A} = the Borel σ -algebra of \underline{C} . Let $\underline{B} = \int_{\mathbb{R}} B \times [0, 1] / B$ is a Borel subset of [0, 1]. Let λ be the product Lebesgue measure on ($(\underline{\zeta})$, \underline{A}). Let λ^* denote the outer Lebesgue measure. Clearly λ is $(\underline{\mathbb{B}}, \underline{\mathbb{A}})$ nonatomic. Let Dbe a set with $\lambda^*(D) = 1$; further let for every $x \in [0, 1]$, $D^{X} = \{y : (x, y) \in \mathcal{L}\}$ be at most a singleton. (A construction of such a set D is available, for instance, in Lemma 2 of [21]). Define λ_1 as in above. Then λ_1 is (D () B D () A nonatomic. Evidently D () A is countably generated with singleton atoms, D $\overline{}$ is countably generated and D is one-sheeted with respect to D () B (as D^{X} is at most a singleton for each $x \in [0, 1]$) with $\lambda_{\gamma}(D) = 1.$

CHAPTER 2

THE METRIC SPACE OF SUB G-ALGEBRAS OF A PROBABILITY SPACE

1. Introduction

$$d(\underline{B}_{1},\underline{B}_{2}) = \sup_{\underline{B}_{1} \in \underline{B}_{1}} \inf_{\underline{B}_{2} \in \underline{B}_{2}} P(\underline{B}_{1} \wedge \underline{B}_{2}) + \sup_{\underline{B}_{2} \in \underline{B}_{2}} \inf_{\underline{B}_{1} \in \underline{B}_{1}} P(\underline{B}_{1} \wedge \underline{B}_{2}),$$

where \underline{B}_1 and \underline{B}_2 are elements of \underline{S} (\underline{C}), \underline{A} , \underline{P}). In that paper, this metric is used to obtain an elegant sufficient condition for 'equiconvergence of martingales'. (More about it in section 3). When there is no room for confusion we will just write \underline{S} or (\underline{S} , \underline{A}) in place of (\underline{S} (\underline{C}), \underline{A} , \underline{P}), \underline{A}), \underline{A}), \underline{A} , \underline{P}), \underline{P} , \underline{P}), \underline{P} , \underline{P} , \underline{P} , \underline

In this chapter we study the topological properties like completeness, compactness, local compactness, separability, connectedness, total disconnectedness, perfectness and dimension of (§, d). Also we consider an isomorphism problem.

Let X be any metric space with a bounded metric ρ . On 2^X (that is, on the class of all closed subsets of X) one introduces what is known as the Hausdorff metric, ρ_1 , by setting

 $\rho_1(C,D) = \sup_{x \in C} \inf_{y \in D} \rho(x,y) + \sup_{y \in D} \inf_{x \in C} \rho(x,y), C,D \in 2^X.$

Some topological properties of $(2^X, \rho_1)$ can be characterised using similar properties of (X, ρ) . (See sections dealing with $(2^X)_m$ of [11] and [12]).

For any $\Lambda \in \underline{A}$, denote by $[\Lambda]$ the class of all measurable sets that are equivalent to A (i.e. those sets which differ from Λ by a set of measure zero). Let $X = \{[\Lambda]/\Lambda \in \underline{A}\}$ and define $\rho([\Lambda], [B]) = P(\Lambda \Lambda B)$, where $[\Lambda], [B] \in X$. (X, ρ) is a metric space. X is usually referred to as the measure algebra $\underline{A}(P)$ [9, P.167]. Define a map (\underline{A}) from \underline{S} to \underline{S} as follows.

 $(+)(B) = \{[B]/B \in B\}, B \in S. (+)(B) \text{ is a closed subset of } (X, p) \text{ and hence an element of } 2^X.$

For any \underline{B}_1 , $\underline{B}_2 \in \underline{S}$,

 $d(\underline{B}_1, \underline{B}_2) = \rho_1(+)(\underline{B}_1), (+)(\underline{B}_2),$ where ρ_1 is the Hausdorff distance as introduced above. Denoting by $\underline{\underline{S}}$ the range of (+) in $\underline{2}^{\underline{X}}$, we, therefore, have that (+) is an isometry between $\underline{\underline{S}}$ and $\underline{\underline{S}}$. We will use this fact to prove some of our results.

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We recall a few definitions. Let $(\ \bigcap$, $\ A$, P) be a probability space. As $\ A$ is said to be an atom if (i) P(A) > 0 and (ii) $B \in \ A$, $B \subset A \implies P(B) = 0$ or P(B) = P(A). The probability P is said to be nonatomic on $\ A$ if there are no atoms; it is said to be completely atomic on $\ A$ if $\ \bigcap$ is the union of atoms. If P is completely atomic on $\ A$ then $\ \bigcap$ is a disjoint union of countably many atoms. For $\ A \in \ A$, P is said to be nonatomic on $\ A$ if (i) $\ P(A) > 0$ and (ii) no subset of $\ A$ is an atom.

For our study, we need the following facts about conditional atoms (as introduced in chapter 1). Let \underline{B}_1 and \underline{B}_2 be elements of \underline{S} with \underline{B}_1 (\underline{B}_2). Let \underline{A} be a conditional atom for (\underline{B}_1 , \underline{B}_2). Since for any \underline{B}_2^0 ε \underline{B}_2 , there exists \underline{B}_1^0 ε \underline{B}_1 with \underline{B}_2^0 (\underline{C}_1) \underline{A}_2 we have

 $P(B_2^{\circ} \triangle B_1^{\circ}) = P[(B_2^{\circ} \triangle B_1^{\circ}) (\overline{\ }) A] + P[(B_2^{\circ} \triangle B_1^{\circ}) (\overline{\ }) A^{\circ}]$ $\leq P(A^{\circ}).$

Therefore, for any $B_2^{\circ} \in \underline{B}_2 \quad \inf_{B_1 \in \underline{B}_1} P(B_2^{\circ} \triangle B_1) \leq P(A^{\circ}).$

Hence $d(\underline{B}_1, \underline{B}_2) \leq P(A^c)$. Let \underline{B}_1 and \underline{B}_2 be elements of \underline{S} with \underline{B}_1 (\underline{B}_2 . Let $\{A_\alpha\}$ be a collection of measurable sets generating a sub σ -algebra whose completion is \underline{B}_2 . Then, it is not hard to see that $A \in \underline{B}_2$ is a conditional atom for $(\underline{B}_1, \underline{B}_2)$ if and only if for every A_α one can find $\underline{B}_\alpha \in \underline{B}_1$ such that $P[(A(), A_\alpha), \Delta(A(), B_\alpha)] = 0$. The proof of

'if' part uses the fact that the class of B ϵ \underline{B}_2 with the property that there is B_1 ϵ \underline{B}_1 such that $P[(A () B) \Delta (A () B_1)] = 0$ is a complete sub σ -algebra containing all A_{α} 's. The proof of 'only if' part is trivial.

In the sequel, (\bigcap), \triangle , P) will stand for a fixed probability space and X for the measure algebra \triangle (P). $\widehat{\underline{S}}$ will stand for the image in 2^X of $\widehat{\underline{S}}$ under the mapping considered above. Given any \underline{A} ε \underline{A} with $\underline{P}(A) > 0$, the symbol \underline{P}_{A} will stand for the probability measure on \underline{A}_{O} ($\overline{}$) \underline{A}_{O} , defined by \underline{P}_{A} (B($\overline{}$)) \underline{A}_{O}) = $\frac{\underline{P}(B(\overline{}) \cdot \underline{A}_{O})}{\underline{P}(A)}$, $\underline{B} \varepsilon$ \underline{A}_{O} . Given a collection $\{\underline{A}_{C}, \ \alpha \varepsilon \ | \ \}$ of \underline{A}_{O} -measurable sets, the symbol σ $\{\underline{A}_{C}, \ \alpha \varepsilon \ | \ \}$ will denote the smallest σ -algebra containing all \underline{A}_{C} 's. A word about the symbol ($\underline{+}$). This has not been reserved exclusively for the mapping considered above. In the latter sections of this chapter we have occasions to use ($\underline{+}$) to denote other mappings as well; however, it will be clear from the context which mapping is being referred to.

2. Completeness of (§, d)

It is known that X is complete and hence 2^X is complete [11, the theorem in P.407]. So, to prove that (\S,d) is complete, it is enough to prove that $\widehat{\S}$ is a closed subset of 2^X .

Lemma 2.1 \underline{S} is a closed subset of 2^{X} .

Proof: Let \underline{B}_{n} be a sequence of elements of \underline{S} converging to an element \underline{F} of 2^{X} . Since \underline{F} is a closed subset of X, $[F] \in \underline{F}$ if and only if $\inf_{[G] \in \underline{F}} P(F \triangle G) = 0$. We have G = 0 to show that $\underline{F} \in \underline{S}$. For this, in the light of the above observation, enough to show that

(i)
$$[F] \in \underline{F} \implies \inf_{[G] \in \underline{F}} P(F^C \triangle G) = 0$$
 and

(ii)
$$[F_1], [F_2], \dots, [F_n], \dots, \varepsilon_{\underline{F}} \Longrightarrow \inf_{[G] \in F} P[((\underline{\ }) F_n) \land G]=0.$$

Let $\epsilon > 0$ be given. Let $[F] \epsilon \underline{F}$. Since $\rho_1(\underline{B}_n, \underline{F}) \to 0$, there exists \underline{B}_n , for some sufficiently large n_0 , such that $\rho_1(\underline{B}_n, \underline{F}) < \frac{\epsilon}{2}$. Choose $[B_n] \epsilon \underline{B}_n$, for which $P(F \Delta B_n) < \epsilon/2$. Then

$$\inf_{[G] \in \underline{F}} P(F^{c} \land G) \leq P(F^{c} \land B_{n_{o}}^{c}) + \inf_{[G] \in \underline{F}} P(B_{n_{o}}^{c} \land G)$$

$$= P(F \land B_{n_{o}}) + \inf_{[G] \in \underline{F}} P(B_{n_{o}}^{c} \land G)$$

$$< \frac{\varepsilon}{2} + \rho_{1} (\underline{F}, \underline{B}_{n_{o}}) < \varepsilon.$$

Now, & being arbitrary, proof of (i) follows.

Let $[F_1]$ and $[F_2]$ belong to \underline{F} . Choose n_0 large enough that $p_1(\underline{F}, \underline{B}_{n_0}) < \frac{\varepsilon}{4}$. This implies that there

exist $[B_1]$ and $[B_2]$ in B_0 such that $P(F_j \triangle B_j) < \frac{\varepsilon}{4}$, for j = 1, 2.

 $\inf_{[G] \in \underline{F}} P[(F_1 (_) F_2) \land G] \leq P[(F_1 (_) F_2) \land (B_1 (_) B_2)] + \inf_{[G] \in \underline{F}} P[(B_1 (_) B_2) \land G]$ $[G] \in \underline{F}$ $\leq P(F_1 \land B_1) + P(F_2 \land B_2) + \rho_1 (\underline{F}, \underline{B}_{n_0})$ $< 3 \frac{\varepsilon}{4} < \varepsilon .$

 ε being arbitrary, we have inf P[(F₁ (_) F₂) Δ G] = 0; [G] ε F

i.e. $[F_1] \in F_2$ ϵF . So, by induction

 $[F_1], [F_2], \dots, [F_n] \in \underline{F} \Rightarrow [\underbrace{n}_{m=1}, F_m] \in \underline{F}.$

Now, let $\{[F_n], n \ge 1\}$ $\subseteq F$. Choose n_0 so large that $P(()_{m=1}^{\infty} F_m - ()_{m=1}^{\infty} F_m) < \epsilon$. Since $[()_{m=1}^{\infty} F_m] \epsilon F$. We have,

inf P[((_) F_m) Δ G] < ϵ . Again, ϵ being arbitrary, the proof of (ii) follows.

Theorem 2.2 (S, d) is a complete metric space.

Proof. The proof follows from lemma 2.1 and the observation made in the beginning of this section.

Remark: Boylan gives a different proof of this theorem in [7].

3. Compactness and local compactness of ($\underline{\underline{S}}$, d)

In this section we show that the following are equivalent.

- (i) P is completely atomic on A.
- (ii) (S, d) is compact.
- (iii) (S, d) is locally compact.

Let $B_0 \in \underline{A}$ be such that P is nonatomic on B_0 . Let \underline{B}_0 be the complete sub σ -algebra of \underline{A} defined by

$$\underline{B}_{o} = \{ A \in \underline{A} / P(A (\overline{\ }) B_{o}) = P(A) \text{ or } P(A (\overline{\ }) B_{o}^{c}) = P(B_{o}^{c}) \}.$$

Then we have the following lemma.

Lemma 3.1 Let ϵ be such that $0 < \epsilon < P(B_0)$. A sequence $[B_n]_{n \ge 1}$, from S, can be found satisfying

- (1) $d(\underline{\underline{B}}_n, \underline{\underline{B}}_0) < \varepsilon$, $n \ge 1$ and
- (ii) $d(\underline{B}_n, \underline{B}_m) \ge \frac{\varepsilon}{4}$, $n, m \ge 1, n \ne m$.

<u>Proof</u>: Since P is nonatomic on B_0 given any measurable subset B of B_0 with 0 < P(B) and given any Y < P(B), one can find a measurable subset B' of B with P(B') = Y. We construct a family

B
11
 B 21
 B 22
 $^{\bullet}$
 $^{\bullet}$
 $^{\bullet}$
 $^{\circ}$
 $^$

of meaurable subsets of $B_{\mathbf{o}}$ as given below.

 B_{11} is a subset of B_0 with $P(B_{11}) = \frac{\varepsilon}{2}$. Having obtained the sets in the first n rows, the sets in the $(n+1)^{th}$ row are obtained as follows.

 $B_{n+1,1}$ is a subset of B_{n1} with $P(B_{n+1,1}) = \frac{P(B_{n1})}{2}$ and $B_{n+1,2} = B_{n1} - B_{n+1,1}$. In general, for $1 \le k \le 2^{n-1}$, the set $B_{n+1,2k-1}$ is a subset of B_{nk} with $P(B_{n+1,2k-1}) = \frac{P(B_{nk})}{2}$ and $B_{n+1,2k} = B_{nk} - B_{n+1,2k-1}$. The rows of (*) form successively finer partitions of B_{11} with $P(B_{n,k}) = \frac{\varepsilon}{2^n}$, $1 \le k \le 2^{n-1}$ and $n \ge 1$.

Now, define

 $\underline{\underline{B}}_{1} = \{\underline{A} \in \underline{\underline{B}}_{0} / P(\underline{A}(\overline{}) B_{11}) = P(\underline{B}_{11}) \text{ or } 0\} \text{ More generally,}$ for $n \ge 1$ $\underline{\underline{B}}_{n} = \{\underline{A} \in \underline{\underline{B}}_{0} / P(\underline{A}(\overline{}) B_{nk}) = P(\underline{B}_{nk}) \text{ or } 0, \text{ for all } 1 \le k \le 2^{n-1}\}.$

Clearly (a) $\{ \underline{B}_n \}_{n \geq 1} \subset \underline{S}$, (b) $\underline{B}_n \subset \underline{B}_0$ for all $n \geq 1$, (c) \underline{B}_n 's are increasing with n and (d) for every n, the set \underline{B}_{11}^c is a conditional atom for $(\underline{B}_n, \underline{B}_0)$. Hence $d(\underline{B}_n, \underline{B}_0) \leq P(\underline{B}_{11}) = \frac{\varepsilon}{2}$. This proves (i).

Let D be the $\underline{\mathbb{B}}_{n+1}$ set defined by $D = \begin{pmatrix} 2^{n-1} \\ k=1 \end{pmatrix}^{B} = n+1, 2k-1$. Observe that for any $B \in \underline{\mathbb{B}}_{n}$, we have $D \triangle B = 0$ D $\triangle (B(\overline{})B_{11})$; and that $(B(\overline{}) B_{11})$ is the union of some of B_{nk} 's. Therefore, it follows that $\inf_{B \in \underline{\mathbb{B}}_{n}} P(D \triangle B) = \frac{\varepsilon}{4}$. Hence $d(\underline{\mathbb{B}}_{n}, \underline{\mathbb{B}}_{n+1}) \geq \frac{\varepsilon}{4}$. As $\underline{\mathbb{B}}_{n}$'s increase with n, for m > n we have $d(\underline{\mathbb{B}}_{n}, \underline{\mathbb{B}}_{n}) \geq d(\underline{\mathbb{B}}_{n}, \underline{\mathbb{B}}_{n+1}) \geq \frac{\varepsilon}{4}$. This proves (ii).

Theorem 3.2 The following are equivalent.

- (i) P is completely atomic on A.
- (ii) (§, d) is compact.
- (iii) (S, d) is locally compact.

Proof. Let P be completely atomic on A. Then X is compact([15, Section 1, p.94] or [2]). So, 2^X is compact ([by the theorem on p.47 and theorem 1, on p.45 of (12)]). So being a closed subset of 2^X, is, therefore compact. Hence S is compact.

The fact that (ii) implies (iii) is immediate.

Suppose P is not completely atomic. So, there exists $B_0 \in \underline{A}$ with $P(B_0) > 0$ and P is nonatomic on B_0 . Then it is clear that, by lemma 3.1, no closed sphere around \underline{B}_0 (as defined before that lemma) is compact. Thus (iii) \Rightarrow (i).

A digression: Theorem 3.2 can be used to make an observation on equiconvergence of martingales. First, the necessary details.

Let $\{\underline{F}_n\}_{n\geq 1}$ be a sequence of elements from \underline{S} , increasing to \underline{F}_{∞} ; i.e. \underline{F}_n 's are increasing and \underline{F}_{∞} is the completion of $(\underline{\ \ \ \ })$ \underline{F}_n . Let $Z \in L_1$ $(\underline{\ \ \ \ })$, \underline{A} , \underline{P}). For any $\underline{B} \in \underline{S}$, let $\underline{P}(Z|\underline{B})$ stand for the conditional expectation of \underline{Z} given \underline{B} . It is part of the folklore that $\{\underline{P}(Z|\underline{F}_n)\}_n \geq 1$ is a martingale and $\underline{P}(Z|\underline{F}_n)$ \underline{L}_1 $\underline{P}(Z|\underline{F}_{\infty})$. (See, theorem \underline{V} \underline{T} 18 of \underline{L}_1). Let \underline{J} be a subset of \underline{L}_1 (\underline{C}), \underline{A} , \underline{P}).

Definition (Boylan, [7]): The martingales

 $\int |P(Z|\underline{F}_n) - P(Z|\underline{F}_\infty)| dP \le \varepsilon$.

The metric d gives an elegant sufficient condition for strong equiconvergence of martingales.

Theorem (Page 554, [7]) Let $J = \{Z/||Z||_{\infty} \le 1\}$. Let $d(\underline{F}_n, \underline{F}_{\infty}) \to 0$. Then the martingales $\{P(Z|\underline{F}_1), P(Z|\underline{F}_2), \dots, P(Z|\underline{F}_n), \dots, /Z \in J\}$ are strongly equiconvergent.

As pointed out by Boylan, it is not true that \underline{F}_n 's increase to \underline{F}_{∞} implies $d(\underline{F}_n, \underline{F}_{\infty}) \to 0$. Also, in this theorem, J can be taken to be any uniformly integrable subset of $L_1(\ C)$, \underline{A} , P).

We want to assert that when \underline{P} is completely atomic on on $\underline{\underline{A}}$, this sufficient condition (viz. $d(\underline{\underline{F}}_n, \underline{\underline{F}}_\infty) \to 0$) is satisfied and hence strong equiconvergence of martingales obtains. For this we need a lemma.

Lemma Let $\underline{\underline{S}}$ be compact. Let $\{\underline{\underline{F}}_n\}_{n \geq 1}$ be a sequence from $\underline{\underline{S}}$ increasing to $\underline{\underline{F}}_{\infty}$. Then $d(\underline{\underline{F}}_n, \underline{\underline{F}}_{\infty}) \rightarrow 0$.

Proof: Since \underline{S} is compact, there is a subsequence $\{n_k\}$ and $\underline{G} \in \underline{S}$ such that $d(\underline{F}_{n_k}, \underline{G}) \to 0$. We will first show $\underline{G} = \underline{F}_{co}$. Let $G \in \underline{G}$. Let $\varepsilon > 0$ be given. Since $d(\underline{F}_{n_k}, \underline{G}) \to 0$, there exists $F \in \underline{F}_{n_k}$ (for some sufficiently large n_k) such that $P(G \land F) < \varepsilon$. So, inf $P(G \land F) < \varepsilon$. Since ε is arbitrary, $F \in \underline{F}_{\infty}$

we have, $\inf_{F \in F_{\infty}} P(G \triangle F) = 0$. Using the fact that F_{∞} is complete, we conclude $G \in F_{\infty}$. Thus $G \subseteq F_{\infty}$. To show $F_{\infty} \subseteq G$, it is enough to show $F_{\infty} \subseteq G$. Let $F \in F_{\infty}$ for some $F_{\infty} \cap F_{\infty} \cap F_{\infty}$ for all sufficiently large $F_{\infty} \cap F_{\infty} \cap F_{\infty}$. Now let $F_{\infty} \cap F_{\infty} \cap F_{\infty}$ for all sufficiently there is $F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty}$. Since $F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty}$ for all sufficiently there is $F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty}$. Therefore $F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty}$ and $F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty}$ are contained in $F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty} \cap F_{\infty}$ are contained in $F_{\infty} \cap F_{\infty} \cap F$

. d($\underline{\underline{F}}_{n_k}$, $\underline{\underline{F}}_{\infty}$) --> 0.

Now, let $\epsilon > 0$ be given. Choose n_k such that $d(\underline{F}_{n_k}, \underline{F}_{\infty}) < \epsilon$. Let $n \ge n_k$. Then

 $\mathrm{d}(\underline{\underline{F}}_n,\ \underline{\underline{F}}_\infty) \leq \mathrm{d}(\underline{\underline{F}}_{n_k},\ \underline{\underline{F}}_\infty) \ (\mathrm{as}\ \underline{\underline{F}}_n \ \widehat{\ }) < \epsilon \ .$

Thus $d(\underline{F}_n, \underline{F}_{\infty}) \rightarrow 0$ and hence the lemma.

Let P be completely atomic on \underline{A} . By theorem 3.2 $\underline{\underline{S}}$ is compact. So, by the lemma above $d(\underline{\underline{F}}_n, \underline{\underline{F}}_\infty)$ -> 0. Hence strong equiconvergence of martingales obtains.

4. Separability of (S, d)

We show that $\underline{\underline{S}}$ is separable if and only if \underline{P} is completely atomic on $\underline{\underline{A}}$.

Consider the two-point space $\{0, 1\}$, equipped with the discrete σ -algebra \underline{D} and the probability measure Q which gives mass $\frac{1}{2}$ to each point. Let N be the set of all natural numbers. Let $\underline{C}_1 = \{0, 1\}^N$, Q^N the product measure on \underline{D}^N and \underline{C} the Q^N -completion of \underline{D}^N .

Lemma 4.1 The space \underline{S} (\underline{C}_1 , \underline{C} , \underline{Q}^N) is not separable.

<u>Proof</u>: Let d_1 denote the distance in $\underline{\underline{S}}$ ($\underline{\Gamma}$), $\underline{\underline{C}}$, \underline{Q}^N).

We shall exhibit an uncountable family $\{\underline{\underline{C}}_{\alpha}\}_{\alpha \in \Gamma}$ of elements of $\underline{\underline{S}}$ ($\underline{\Gamma}$), $\underline{\underline{C}}$, \underline{Q}^N) with $d_1(\underline{\underline{C}}_{\alpha},\underline{\underline{C}}_{\alpha'}) \geq \frac{1}{2}$, $\alpha \neq \alpha'$ and $\alpha,\alpha' \in \Gamma$.

Let

For $m \in \mathbb{N}$, let π_m denote the m^{th} coordinate mapping from C_1 to $\{0,1\}$. Set C_{α} to be the Q^N -completion of $\sigma \{\pi_m^{-1}(\underline{D})/m \text{ is an element of } \alpha\}$. Let $\alpha \neq \alpha'$. Then either α has an element m which is not in α' or vice versa; say, m belongs to α and does not belong to α' . Since the coordinate mappings are independent under Q^N , it is clear that $\pi_m^{-1}(\underline{D})$ is independent of $C_{\alpha'}$. Let $A_0 = \pi_m^{-1}(\{0,0\})$. Then for any B in $C_{\alpha'}$,

$$Q^{N}(A_{O} \triangle B) = Q^{N}(A_{O} - B) + Q^{N}(B - A_{O})
 = Q^{N}(A_{O}) Q^{N}(B^{C}) + Q^{N}(B) Q^{N}(A_{O}^{C})
 = \frac{1}{2}[Q^{N}(B^{C}) + Q^{N}(B)] \text{ (since } Q^{N}(A_{O}) = \frac{1}{2})
 = \frac{1}{2}.$$

So, $\sup_{A \in \underline{C}_{\alpha}} \inf_{B \in \underline{C}_{\alpha'}} \mathbf{Q}^{N}(A \triangle B) \geq \frac{1}{2}$. This implies that $d_{1}(\underline{C}_{\alpha}, \underline{C}_{\alpha'}) \geq \frac{1}{2}$ if $\alpha \neq \alpha'$.

Evidently | is an uncountable set.

Lemma 4.2 Let P be nonatomic on A. Then S is not separable.

<u>Proof:</u> Since P is nonatomic on \underline{A} , we can find a countably generated sub σ -algebra \underline{A} of \underline{A} such that P on \underline{A} is nonatomic. Enough to show that \underline{S} ($\underline{\Gamma}$), completion of \underline{A} , P) is not separable.

Now, by Halmos-Von Neumann theorem [9, theorem C p.173], the measure algebra $\underline{A}_{0}(P)$ is isomorphic to the measure algebra $\underline{D}^{N}(Q^{N})$ (as introduced in the beginning of this section). This establishes an isometry between $\underline{S}(\underline{C})$, completion of $\underline{A}_{0},P)$ and $\underline{S}(\underline{C})_{1},\underline{C}_{1},Q^{N})$. (For details see section 9). Since $\underline{S}(\underline{C})_{1},\underline{C}_{1},Q^{N}$ is not separable by lemma 4.1, the desired conclusion follows. Theorem 4.3 \underline{S} is separable if and only if \underline{P} is completely atomic on \underline{A}_{0} .

Proof: Let P be completely atomic on A. Then by theorem 3.2 (S, d) is compact and hence separable.

Let now there exist $B_0 \in \underline{A}$ such that P is nonatomic on B_0 . From lemma 4.2 it follows that $\underline{S}(B_0, B_0) = (1) \cdot \underline{A}$, P_{B_0} is not separable. Define a mapping (+) from $\underline{S}(\underline{B}_0, B_0) = (1) \cdot \underline{A}$, P_{B_0} to \underline{S} by setting for $\underline{D} \in \underline{S}(B_0, B_0) = (1) \cdot \underline{A}$, \underline{P}_{B_0}

 $(+)(\underline{D}) = P\text{-completion of } \{A/A \in \underline{D} \text{ or } \underline{C}\} - A \in \underline{D}\}.$ It is not difficult to check that (+) is a homeomorphism between $\underline{S}(B_0, B_0, C) \triangleq P_{B_0}$ and $(+)[\underline{S}(B_0, B_0, C) \triangleq P_{B_0}].$ This implies that a subset of \underline{S} is not separable. Hence \underline{S} is not separable.

5. Connectedness and total disconnectedness of (\S, d) .

Lemma 5.1 Let P be such that it has at most one atom. Let A_0 stand for the atom, if there is one; for the empty set, otherwise. Then, given any $B \in S$ a continuous function f can be defined from the interval $[0, 1 - P(A_0)]$ to S satisfying

(ii)
$$f(1 - P(A)) = A$$
.

⁽i) $f(0) = \underline{B}$ and

<u>Proof</u>: From the hypothesis of the lemma we have that P is nonatomic on A_0^C . Then, it is well-known that we can find a collection $\{B_t\}_{t \in [0, 1-P(A_0)]}$ of measurable sets satisfying

(cc)
$$B_{t_1} \subset B_{t_2}$$
 if $t_1 < t_2$,

$$(\beta)$$
 $P(B_t) = t$ and

$$(\gamma) \quad B_{1-P(A_0)} = A_0^{c}$$

where t, t_1 , t_2 are from [0, 1 - P(A₀)]. Define complete sub σ -algebras \underline{B}_t of \underline{A} by

 $\underline{B}_{t} = \{A \in \underline{A} / P(A (\overline{\ }) B_{t}^{c}) = P(B_{t}^{c}) \text{ or } 0\}, \text{ t } \in [0, 1 - P(A_{0})].$ If $t_{1} < t_{2}$, we have $B_{t_{1}} \subset B_{t_{2}}$ and so $\underline{B}_{t_{1}} \subset \underline{B}_{t_{2}}$. Now, $\underline{B}_{t_{1}} (\underline{\ }) B_{t_{2}}^{c} \text{ is a conditional atom for } (\underline{B}_{t_{1}}, \underline{B}_{t_{2}}). \text{ To see}$ this, let $A \in \underline{B}_{t_{2}}$. If $P(A (\overline{\ }) B_{t_{2}}^{c}) = 0$, consider the $\underline{B}_{t_{1}}$ set A_{1} defined by $A_{1} = A (\overline{\ }) B_{t_{1}}$. Then $P[A (\overline{\ }) (B_{t_{1}} (\underline{\ }) B_{t_{2}}^{c}) \triangle A_{1} (\overline{\ }) (B_{t_{1}} (\underline{\ }) B_{t_{2}}^{c})] = 0. \text{ If}$ $P(A (\overline{\ }) B_{t_{2}}^{c}) = P(B_{t_{2}}^{c}), \text{ consider the } \underline{B}_{t_{1}} \text{ set } A_{1} \text{ defined by}$ $A_{1} = A (\underline{\ }) (B_{t_{1}}^{c} - B_{t_{2}}^{c}). \text{ Then}$ $P[A (\overline{\ }) (B_{t_{1}} (\underline{\ }) B_{t_{2}}^{c}) \triangle A_{1} (\overline{\ }) (B_{t_{1}} (\underline{\ }) B_{t_{2}}^{c})] = 0.$

Now define complete sub σ -algebras $\underline{\mathcal{C}}_t$ of $\underline{\underline{\mathbf{A}}}$ by setting

 $\underline{\underline{C}}_t$ = the completion of $\sigma\{\underline{\underline{B}}(\underline{\underline{}},\underline{\underline{B}}_t\}$, to $[0,1-P(A_0)]$.

For $t_1 < t_2$, the set $B_{t_1}(\underline{\underline{}},\underline{\underline{B}}_t)$ is a conditional atom for $(\underline{\underline{C}}_{t_1},\underline{\underline{C}}_{t_2})$. To verify this, in view of an observation made in section 1, we have to only check that for any $\underline{\underline{A}}$ in $\underline{\underline{B}}_{t_2}$ or $\underline{\underline{B}}$ there exists $\underline{\underline{A}}_1$ in $\underline{\underline{C}}_{t_1}$ such that $P[\underline{\underline{A}}(\underline{\underline{}},\underline{\underline{C}}_{t_2})]$ and $\underline{\underline{A}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ by $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ and $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ by $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ and $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ by $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ in $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ by $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}}_t)$ in $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}},\underline{\underline{C}}_t)$ by $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}},\underline{\underline{C}}_t)$ in $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}},\underline{\underline{C}}_t)$ in $\underline{\underline{C}}_1(\underline{\underline{}},\underline{\underline{C}},$

 $d(\underline{\mathcal{G}}_{t_1}, \underline{\mathcal{G}}_{t_2}) \leq P[(B_{t_1}, \underline{\mathcal{G}}_{t_2})^c] = P(B_{t_2}, \underline{\mathcal{G}}_{t_1}) = t_2 - t_1.$ Hence the function f defined on $[0, 1 - P(A_0)]$ by $f(t) = \underline{\mathcal{G}}_t$ is continuous.

$$B_{O} = \{A \in \underline{A} / P(A(\overline{\ }) B_{O}^{C}) = P(B_{O}^{C}) \text{ or } 0 \}$$

$$= \{A \in \underline{A} / P(A) = 1 \text{ or } 0\} \text{ and}$$

$$B_{I}-P(A_{O}) = \{A \in \underline{A} / P(A(\overline{\ }) B_{I}^{C}-P(A_{O})) = P(B_{I}^{C}-P(A_{O})) \text{ or } 0\}$$

$$= \{A \in \underline{A} / P(A(\overline{\ }) A_{O}) = P(A_{O}) \text{ or } 0\}$$

$$= \{A \in \underline{A} / P(A(\overline{\ }) A_{O}) = P(A_{O}) \text{ or } 0\}$$

$$= A_{O}, \text{ since } A_{O} \text{ is } \text{the } \text{only atom.}$$
So, $f(O) = B_{O} \text{ and } f(I - P(A_{O})) = A_{O}.$

Theorem 5.2 The following are equivalent

- (i) P has at most one atom.
- (ii) (S, d) is arcwise connected.
- (iii) (S, d) is connected.

Proof: Lemma 5.1 gives the proof of (i) ⇒ (ii). That (ii) ⇒ (iii) is trivial.

Let P have at least two atoms, say A_1 and A_2 . Consider the probability space $(A_1(_) A_2, (A_1(_) A_2) (\overline{}) A_1 (\underline{}) A_2$.

Clearly, $\underline{S}_1 = \underline{S}(A_1 (_) A_2, (A_1 (_) A_2) (\overline) \underline{A}, P_{A_1} (_) A_2)$ contains exactly two elements. Define the map (+) from \underline{S} to \underline{S}_1 by setting, for $\underline{B} \in \underline{S}$, $(+)(\underline{B}) = (A_1 (_) A_2) (\overline{)} \underline{B}$.

(+) is continuous and onto \underline{S}_1 . Since \underline{S}_1 is not connected, \underline{S}_1 is not connected. Thus (iii) \Longrightarrow (i). Hence the theorem.

Regarding total disconnectedness we have

Theorem 5.3 (\S , d) is totally disconnected if and only if P is completely atomic.

Proof: P is completely atomic implies that X is compact and totally disconnected ([2], theorem 6.1). This implies that 2^{X} is totally disconnected (by the theorem on p.47 of [12] and proposition 4.13.2 of [17]). So, \widehat{S} and hence S is totally disconnected.

Now, let P be nonatomic on $B_0 \in \underline{A}$. We shall show that there is a connected subset of (\underline{S}, d) containing at least two elements. This would then show that (\underline{S}, d) cannot be totally disconnected.

 $(+)(\underline{D}) = P$ -completion of $\{ \Lambda \in \underline{A} / \Lambda \in \underline{D} \text{ or } \underline{C} \}$

The mapping (+) can be verified to be a continuous one. Therefore $(+)(\underline{S}_{\underline{1}})$ is a connected subset of \underline{S} , being the continuous image of a connected set.

 $(+) (\{\emptyset, B_o\}) = P\text{-completion of } \{\emptyset, B_o, B_o^c, C\} \} \text{ and}$ $(+) (B_o(\overline{\ }) \underline{A}) = \{A \in \underline{A} / P(A(\overline{\ }) B_o) = P(A) \text{ or } P(A(\overline{\ }) B_o^c) = P(B_o^c) \}.$

The two σ -algebras above are obviously different. Thus (+) (\underline{S}_1) contains at least two points. Hence the theorem.

In the remainder of this section we study the components of (\S , d).

Let A_0 be the union of all atoms. (A may be empty). Consider (A_0, A_0) (A_1, P_{A_0}) . (In case A_0 is empty take (A_0, A_0) (A_1, P_{A_0}) to be a set with one element). Define a mapping (A_1) from A_1 to A_2 (A_1) (A_1)

Theorem 5.4 The components of \underline{S} are precisely the collection $(+)^{-1}(\underline{D})$: $\underline{D} \in \underline{S}_1$, where (+) is the mapping as defined above.

Proof. We merely sketch the proof, omitting the details to the reader.

Since (+) is continuous and \underline{S}_1 is totally disconnected, the image under (+) of any component of \underline{S} is a singleton. So, we have to show only that for every $\underline{D} \in \underline{S}_1$, the set $(+)^{-1}(\underline{D})$ is connected in \underline{S} . Fix $\underline{D}_0 \in \underline{S}_1$. We shall show that $(+)^{-1}(\underline{D}_0)$ is arcwise connected. Let $\underline{F}_1 = \{A \in \underline{A} / A(\overline{}) \ A_0 \in \underline{D}_0\}$ and let \underline{F}_0 be an arbitrary element of $(+)^{-1}(\underline{D}_0)$. Consider $\{\underline{B}_t\}_t \in [0,1-P(A_0)]$ as defined in lemma 5.1 and set $\underline{C}_t = \text{completion of } \sigma(\underline{B}_t(-),\underline{F}_0)$. It can be verified that for each \underline{t} , $\underline{C}_t \in (+)^{-1}(\underline{D}_0)$. Define \underline{f} from $[0, 1-P(A_0)]$ to $(+)^{-1}(\underline{D}_0)$ by $\underline{f}(t) = \underline{C}_t$. Then,

f is a continuous map with $f(0) = \underline{F}_0$ and $f(1-P(A_0)) = \underline{F}_1$. Since \underline{F}_0 is arbitrary, the result follows.

6. Perfectness of (S, d)

We show that (\underline{S} , d) is perfect if and only if the range of P, i.e. the set $P(\Lambda)/\Lambda \in \underline{A}$, is an infinite set.

Lemma 6.1 Let \underline{B}_1 and \underline{B}_2 be elements of \underline{S} such that for some $\underline{A} \in \underline{A}$, \underline{A}^c () $\underline{B}_1 = \underline{A}^c$ () \underline{B}_2 . Then $\underline{d}(\underline{B}_1, \underline{B}_2) \leq 2P(\underline{A})$.

<u>Proof:</u> The proof of this is similar to the proof of a result in the introduction.

Lemma 6.2 Let the range of P be an infinite set. Then given $\epsilon > 0$, one can find $A \in \underline{A}$ with $0 < P(A) < \frac{\epsilon}{2}$.

Proof: Suppose there is $A \in \underline{A}$ such that P restricted, to A is nonatomic; then the proof is clear. Let now P be completely atomic. Since the range of P is an infinite set we can find $\{A_n\}_n \geq 1$ of atoms which are pairwise disjoint. The fact that $\sum_{n=1}^{\infty} P(A_n) = 1$ gives the desired conclusion.

Theorem 6.3 (S, d) is perfect if and only if the range of P is an infinite set.

<u>Proof</u>: Let the range of P be a finite set. Then \underline{S} contains only finitely many elements and hence is not perfect.

Let the range of P be an infinite set. Let $\underline{\underline{B}}_0 \in \underline{\underline{S}}$ and $\epsilon > 0$ be given. We have to exhibit a $\underline{\underline{B}}_1 \in \underline{\underline{S}}$ such that $0 < d(\underline{\underline{B}}_0, \underline{\underline{B}}_1) < \epsilon$.

Case (i) P on \underline{B}_0 has a nonatomic part. Then choose $\Lambda_0 \in \underline{B}_0$ such that $0 < P(\underline{A}_0) < \frac{\varepsilon}{2}$ and no \underline{B}_0 measurable subset of \underline{A}_0 is atom for P. Let \underline{B}_1 be the completion of $\sigma(\Lambda_0, \Lambda_0^C(\overline{}), \underline{B}_0)$. Because of nonatomicity of P on \underline{A}_0 we have $0 < d(\underline{B}_0, \underline{B}_1)$ and by lemma 6.1, $d(\underline{B}_0, \underline{B}_1) < \varepsilon$.

Case (ii) P on \underline{B}_{G} is completely atomic and has only finitely many atoms.

Let A_1 , A_2 , ..., A_m be all the atoms. By lemma 6.2, we can get $A_0 \in \underline{A}$ with $0 < P(A_0) < \min \{P(A_1), P(A_2), \ldots, P(A_m), \frac{\varepsilon}{2}\}$. Take the required \underline{B}_1 to be the completion of $\sigma(A_0, A_0^C(\underline{C}), \underline{B}_0)$. By the choice of A_0 , $d(\underline{B}_0, \underline{B}_1) > 0$ and by lemma 6.1, $d(\underline{B}_0, \underline{B}_1) < \varepsilon$.

Case (iii) P on \underline{B}_0 is completely atomic and has infinitely many atoms.

Choose two atoms A_1 and A_2 with $0 < P(A_1) + P(A_2) < \frac{1}{2}$. (This can be done since there are infinitely many atoms). Let $A_0 = A_1$ () A_2 and define B_1 by $B_1 = \text{Completion of}$ of A_0 , A_0^c () A_0 and by the choice of A_0 , A_0^c () A_0 and by lemma 6.1 A_0 (A_0) A_0) A_0 (A_0) A_0 0 (A_0) A_0 0 (A_0) A_0 0 (A_0 0) A_0 0 (A_0

This completes the proof of 'if' part. Hence the theorem.

7. Dimension of (S, d)

For notions in dimension theory we refer to Nagata [18]. By dimension we mean covering dimension in the sense of Nagata [See 18, p.9]. We need the following theorem from [2].

Theorem 7.1 Let P be nonatomic on \underline{A} . Then dimension of \underline{A} (P) is infinity.

For a proof see theorem 9.2 of [2].

Regarding the dimension of (\underline{S}, d) we have the following theorem.

Theorem 7.2 The dimension of (S, d) is either zero or infinity. If P is completely atomic on A the dimension is zero; otherwise it is infinity.

<u>Proof.</u> Let P be completely atomic on \underline{A} . Then, by theorem 5.3, (\underline{S} , d) is totally disconnected; hence the dimension of (\underline{S} , d) is zero.

Let now there exist $B_0 \in \underline{A}$, such that P is nonatomic on B_0 . Without loss in generality we can take that $0 < P(B_0) < \frac{1}{2}$. Consider the measure algebra $(B_0()\underline{A})(P_B)$. By theorem 7.1 the dimension of $(B_0()\underline{A})(P_b)$ is infinity. We shall establish a

homeomorphism (+) from (B₀ () (P_B) into $\underline{\underline{S}}$ such that (+)[(B₀ () $\underline{\underline{A}}$) (P_B)] is a closed subset of $\underline{\underline{S}}$. Then it would follow that the dimension of $\underline{\underline{S}}$, being greater than or equal to that of (+)[(B₀ () $\underline{\underline{A}}$) (P_B)], is infinity.

Define the map (+) from $(B_0 \cap \underline{A})(P_{B_0})$ to \underline{S} by setting, for $[B_1] \in (B_0 \cap \underline{A})(P_{B_0})$ $(+)([B_1]) =$ the P-completion of $\{\emptyset, B_1, \Omega - B_1, \Omega \}$.

The fact that (+) is well defined is easily verified. For two elements B_1 and B_2 of $\underline{A}(\overline{B}) = 0$ we have $(A_1) = 0$ $A_2 = 0$ and $(A_1) = 0$ $A_2 = 0$ $A_3 = 0$ $A_4 = 0$

Thus d($(+)(B_1)$, $(+)(B_2)$) = min[P(B_1), P($B_1 \triangle B_2$)] + min[P(B_2), P($B_1 \triangle B_2$)].

Hence $P(B_1 \triangle B_2) \le d((+)(B_1), (+)(B_2)) \le 2P(B_1 \triangle B_2) - (**)$

Evidently, (+) is a 1-1 map. Using (**) we conclude that (+) is a homeomorphism between (B₀ (-) \triangle)(P_{B₀}) and (+)[(B₀ (-) \triangle)(P_{B₀})] is a

complete subset of \S . Since \S is complete, we have that $(+)[(B_0(\overline{\ }) \triangle)(P_B)]$ is closed.

In view of the observation made in the earlier part of the proof, we have that if P is not completely atomic the dimension of (S, d) is infinity.

Hence the theorem.

8. Some additional results on (S, d)

In this section we show that the following two subsets of \underline{S} , namely

 $\underline{\underline{S}}_{c.a.} = \{\underline{B} \in \underline{\underline{S}}./\underline{P} \text{ is completely atomic on } \underline{B} \}$ and $\underline{\underline{S}}_{s} = \{\underline{B} \in \underline{\underline{S}}/\underline{B} \text{ is the completion of a countably generated sub } \sigma \text{-algebra } \}$ are closed.

Let P be completely atomic on \underline{B}_0 with infinitely many atoms. Let $\{A_n\}_{n\geq 1}$ be all the atoms. Let $\epsilon>0$ be given. Choose n_0 such that $\sum_{m\geq n} P(A_m) < \epsilon$. Define $\sum_{m\geq n} = \{ b \in \mathbb{Z} \text{ and } P \text{ on } B \text{ has only finitely many atoms} \}$ is dense in $\sum_{n\geq n} a_n = \{ b \in \mathbb{Z} \text{ a. } e \}$.

Lemma 8.1 Let $\underline{\underline{B}}_0 \in \underline{\underline{S}}$. Let $\{\underline{\underline{B}}_n\}_{n \geq 1} \subset \underline{\underline{S}}_{c.a.}^i$ be such that for each n (i) $\underline{\underline{B}}_n \subset \underline{\underline{B}}_0$ and (ii) $d(\underline{\underline{B}}_n, \underline{\underline{B}}_0) \to 0$ as $n \to \infty$. Then $\underline{\underline{B}}_0 \in \underline{\underline{S}}_{c.a.}$.

<u>Proof:</u> Suppose there exists $B_0 \in \underline{B}_0$ such that $P(B_0) > 0$ and no \underline{B}_0 measurable subset of B_0 is an atom.

Fix $n(\geq 1)$ and consider \underline{B}_n . Let $\underline{A}_1, \underline{A}_2, \dots, \underline{A}_k$ be all the atoms of \underline{B}_n . For each $j(1 \leq j \leq k)$, choose $\underline{B}_j \in \underline{J}$ such that $\underline{B}_j = \emptyset$ if $P(\underline{A}_j \cap B_0) = 0$; $\underline{B}_j \cap A_j \cap B_0$ and $P(\underline{B}_j) = \frac{P(\underline{A}_j \cap B_0)}{2}$ if $P(\underline{A}_j \cap B_0) > 0$. (Such a choice is possible since $\underline{B}_n \cap B_0 = 0$). Let $\underline{B}_n = 0$ and $\underline{B}_n \cap B_0 = 0$. So, $\underline{A}_n \cap B_0 = 0$ and $\underline{A}_n \cap B_0 = 0$. So, $\underline{A}_n \cap B_0 = 0$ as a representation, the since $\underline{A}_n \cap B_0 = 0$ as $\underline{A}_n \cap B_0 = 0$, by hypothesis. Hence the lemma.

Theorem 8.2 $\subseteq_{c.a.}$ is a closed subset of \subseteq .

Proof: Let $(\subseteq_n]_n \ge 1$ $\subseteq_{c.a.}$ be such that for some $\subseteq_{c.a.}$ be such that for some $\subseteq_{c.a.}$ be such that $\subseteq_{c.a.}$ be such that for some $\subseteq_{c.a.}$ be such that $\subseteq_{c.a.}$

assume without loss in generality that $\{\underline{c}_n\}_{n \geq 1} \subset \underline{S}_{c,a}$.

Fix $n(\geq 1)$. Let A_1, A_2, \dots, A_k be all the atoms of \underline{C}_n . For each subset J of $\{1,2,\dots,k\}$ choose a set B_J^n of \underline{B}_0 such that $P(B_J^n \land \bigcup_{j \in J} A_j) < 2d(\underline{C}_n, \underline{B}_0)$. Let $\underline{B}_n = P$ -completion of $\sigma(B_J^n : J)$ is a subset of $\{1,2,\dots,k\}$. Now, for any B in \underline{B}_0 , there exists some subset J of $\{1,2,\dots,k\}$ such that $P(B \land \bigcup_{j \in J} A_j) < 2d(\underline{C}_n, \underline{B}_0)$. So, for any B in \underline{B}_0 , there exists B_J^n in \underline{B}_n such that $P(B \land B_J^n) < 4d(\underline{C}_n, \underline{B}_0)$. Hence $P(B \land B_J^n) < 4d(\underline{C}_n, \underline{B}_0)$. Also, we have $\underline{B}_n \subset \underline{B}_0$. Since $P(B \land B_0^n) = P(B \land B_0^n) = P$

Theorem 8.3 S is a closed set.

- Proof. (a) Let $\underline{D} \in \underline{S}_{S}$? Then the metric space $\underline{D}(P)$ is separable. Let \underline{C} (\underline{D}). Then the metric space $\underline{C}(P)$, being a subset of a separable metric space, is separable. So $\underline{C} \in \underline{S}_{S}$.
- $(\beta) \quad \text{Let } \{\underline{\underline{\mathbb{C}}}_n\} \subset \underline{\underline{\mathbb{S}}}_s \quad \text{Denote by } \ \, \underline{\underline{\mathbb{C}}}_n \quad \text{the completion}$ of $\sigma\{(\underline{\underline{\mathbb{C}}}_n)\underline{\underline{\mathbb{C}}}_n\}$. It is easily verified that $\underbrace{V}_n \underline{\underline{\mathbb{C}}}_n \in \underline{\underline{\mathbb{S}}}_s$.

(y) Let $\{\underline{C}_n\}$ \subseteq \underline{S}_s be such that for some $\underline{B}_0 \in \underline{S}$ d(\underline{C}_n , \underline{B}_0) \rightarrow 0 as $n \rightarrow \infty$. Let $\underline{B}_0 \in \underline{B}_0$. Then inf $\underline{P}(\underline{B}_0 \land \underline{D}) \leq \inf_{\substack{C \in C \\ n = n}} \underline{P}(\underline{B}_0 \land \underline{D}) \leq \inf_{\substack{C \in C \\ n = n}} \underline{P}(\underline{B}_0 \land \underline{C}_n)$ for every $n \in \underline{S}_0$,

 $\inf_{\substack{D \in V \subseteq \\ n=n}} P(B \land D) = 0. \text{ This gives that } \underline{B}_{O} \subset V \subseteq_{n}^{V}.$

Now from (α), (β) and (γ) the result follows.

9. Isomorphism Problem

Let us start with an observation. Consider two metric spaces Y_1 and Y_2 equipped with bounded metrics. Let T be an isometry from Y_1 onto Y_2 . Let the mapping Y_1 be defined from $2^{\frac{Y_1}{1}}$ to $2^{\frac{Y_2}{2}}$ by

(+)(C) = $\{T(y_1): y_1 \in Y_1\}$, $C \in 2^{Y_1}$. Then, it is easily seen that (+) is onto 2^{Y_2} and is an isometry.

Let (C_1, A_1, P_1) and (C_2, A_2, P_2) be two probability spaces. Let T be an isomorphism between the measure algebras (P_1) and (P_2) (see page 167 of [9]). i.e. T is a one to one transformation from (P_1) onto (P_2) such that $(E_1 - [C]) = (E_1) - (E_2)$ and

 $T((\bigcup_{n=1}^{\infty} [A_{n1}]) = (\bigcup_{n=1}^{\infty} T([A_{n1}])$ whenever B, C and A_{n1} 's are

elements of \underline{A}_1 ; moreover \underline{T} 'preserves measure' in the sense whenever $\underline{A}_{12} \in \underline{T}[\underline{A}_{11}]$ we have $\underline{P}_1(\underline{A}_{11}) = \underline{P}_2(\underline{A}_{12})$. Evidently, \underline{T} is an isometry from $\underline{A}_1(\underline{P}_1)$ onto $\underline{A}_2(\underline{P}_2)$. Use the observation in the first paragraph of this section to define an isometry $\underline{(+)}$ from $\underline{2}^{\underline{A}_1(\underline{P}_1)}$ onto $\underline{2}^{\underline{A}_2(\underline{P}_2)}$. Restricting $\underline{(+)}$ to $\underline{\underline{S}}(\underline{C}_{11},\underline{A}_{11},\underline{P}_{11})$ we find it is mapped onto $\underline{\underline{S}}(\underline{C}_{12},\underline{A}_{21},\underline{P}_{21})$. Using the natural mapping between $\underline{\underline{S}}$ and $\underline{\underline{S}}$ mentioned in the introduction (page 35), we can look upon $\underline{(+)}$ as an isometry from $\underline{\underline{S}}(\underline{C}_{12},\underline{A}_{11},\underline{P}_{11})$ onto $\underline{\underline{S}}(\underline{C}_{12},\underline{A}_{21},\underline{P}_{21})$ and we do so. It is routine to verify that $\underline{(+)}$ is a complete lattice isomorphism as well.

i.e. for $\{\underline{B}_{\alpha}\}$ \subseteq \underline{S} $(\underline{\Gamma})_1$, \underline{A}_1 , \underline{P}_1),

$$(+)(V_{\alpha}B_{\alpha}) = V_{\alpha}(+)(B_{\alpha}) \text{ and } (+)(\Lambda_{\alpha}B_{\alpha}) = \Lambda_{\alpha}(+)(B_{\alpha}).$$

(Here, 'V' of a family of sub σ -algebras denotes the smallest completed sub σ -algebra containing them and ' $_{\Lambda}$ ' of a family denotes their intersection). We call this mapping ($_{\downarrow}$) as being induced by T. Now, we ask ''Is every isometry ($_{\downarrow}$) from $\underline{S}(\Omega_1,\underline{A}_1,P_1)$ onto $\underline{S}(\Omega_2,\underline{A}_2,P_2)$ which preserves lattice operations induced by an isomorphism T between the measure algebras $\underline{A}_1(P_1)$ and $\underline{A}_2(P_2)$?''. The remainder of this section is intended to provide an affirmative answer to this question.

Let $\{A_{\alpha}, \, \alpha \in \Gamma\}$ be a collection of measurable sets of a probability space ($\{C\}$, A, P). In this section, for ease in presentation, we will let $\sigma\{A_{\alpha}, \alpha \in \Gamma\}$ stand for the P-completion of the σ -algebra generated by $\{A_{\alpha}, \alpha \in \Gamma\}$.

Before proceeding with the proof two remarks are in order.

In any $\underline{S}(\underline{\Gamma})$, \underline{A} , \underline{P}) the distance between any $\underline{B} \in \underline{S}$ and $\sigma(0,\underline{\Gamma})$ is attained; that is, there is $\underline{B}_0 \in \underline{B}$ with $\underline{P}(\underline{B}_0) = \underline{d}(\underline{B}, \sigma(0,\underline{\Gamma}))$. (This follows because (*) $\underline{d}(\underline{B}, \sigma(0,\underline{\Gamma})) = \sup\{\underline{P}(\underline{B}) \mid \underline{B} \in \underline{B} \text{ and } \underline{P}(\underline{B}) \leq \underline{1}\}$ and $\underline{P}(\underline{B}) \mid \underline{B} \in \underline{B} \text{ and } \underline{P}(\underline{B}) \leq \underline{1}\}$ is a closed set).

Also if $\underline{B} = \sigma(A)$ with $P(A) < \frac{1}{2}$ then $d(\underline{B}, \sigma(\emptyset, \Omega)) = P(A)$. (**)

In the sequel (\underline{S}_i, d_i) (or simply \underline{S}_i), i = 1, 2 will stand respectively for $S(\underline{C}_i, \underline{A}_i, P_i)$, i = 1, 2. $A_{01}, A_{11}, A_{21}, \ldots$, A_{n1}, \ldots will be sets from \underline{A}_1 and $A_{02}, A_{12}, A_{22}, \ldots$, A_{n2}, \ldots will be sets from \underline{A}_2 .

Lemma 9.1 Let $\{\underline{B}_{\alpha}\}$ (\underline{S}_{1}) . Then

$$(+)(V_{\alpha}) = V_{\alpha} + (+)(B_{\alpha})$$
 and $(+)(\Lambda_{\alpha}) = \Lambda_{\alpha} + (+)(B_{\alpha})$.

Proof: Since (+) preserves lattice operations and since $\underline{\underline{B}}_{\alpha} \subset \underline{\underline{V}}_{\alpha} \underline{\underline{B}}_{\alpha}$, we have $(+)(\underline{\underline{B}}_{\alpha}) \subset (+)(\underline{\underline{V}}_{\alpha} \underline{\underline{B}}_{\alpha})$ for each α .

$$\infty$$
, ∇ $(+)$ $(\underline{\mathbb{P}}_{\alpha})$ $(-)$ $(\nabla$ $\underline{\mathbb{P}}_{\alpha})$.

Let $\{\underline{\mathcal{G}}_{\alpha}\}$ \subset $\underline{\mathbb{S}}_{2}$. Working with $(+)^{-1}$ instead of (+) we have $(+)^{-1}(\underline{\mathcal{G}}_{\alpha})$ $(-)^{-1}(\underline{\mathcal{G}}_{\alpha})$. Take $\underline{\mathcal{G}}_{\alpha} = (+)(\underline{\mathbb{B}}_{\alpha})$.

$$\Rightarrow (+)(\bigvee_{\alpha} \mathbb{B}_{\alpha}) \subset \bigvee_{\alpha} (+) (\mathbb{B}_{\alpha}).$$

$$\cdot \cdot \cdot (+) (\begin{smallmatrix} v \\ \alpha \end{smallmatrix}) = \begin{smallmatrix} v \\ \alpha \end{smallmatrix} (+) (\begin{smallmatrix} B \\ \alpha \end{smallmatrix}) .$$

The other part is proved similarly. Hence the lemma.

Lemma 9.2 (i) (+) $(\sigma\{ (1) \}) = \sigma\{ (1) \}$.

(ii) For any $A_{11} \in \underline{A}_1$ with $P_1(A_{11}) < \frac{1}{2}$ there is $A_{12} \in \underline{A}_2$ with $P_1(A_{11}) = P_2(A_{12})$ such that $(+)(\sigma(A_{11})) = \sigma(A_{12})$.

Proof: Observe that $(+)^{-1}(\sigma(\{ \subseteq \}) \supseteq \sigma(\{ \subseteq \}))$.

So, (+) $(\sigma(\Omega_1))(\sigma(\Omega_2))$ (as (+) preserves lattice operations); hence (+) $(\sigma(\Omega_1)) = \sigma(\Omega_2)$. This proves (i).

Let $A_{11} \in \underline{A}_1$ with $P_1(A_{11}) < \frac{1}{2}$. Then $d_1(\sigma \{ \subseteq A_{11} \}, \sigma \{A_{11} \}) = P_1(A_{11})$ by (**). So, $d_2((+)(\sigma \{ \subseteq A_{11} \}), (+)(\sigma \{A_{11} \})) = P_1(A_{11})(as (+) preserves metric).$

i.e. $P_1(A_{11}) = d_2$ (of Ω_2 , (+)(of A_{11})). Using (*) we get $A_{12} \in (+)(o(A_{11}))$ with $P_2(A_{12}) = d_2(o(A_{12}), (+)(o(A_{11})))$. Then we have $P_2(A_{12}) = P_1(A_{11})$. It remains to prove $(+)(o(A_{11})) = o(A_{12})$. Clearly, $(+)(o(A_{11})) \supseteq o(A_{12})$. If the equality is not true, $(+)(o(A_{11}))$ will contain at least two distinct elements of S_2 , both of them different from $o(A_{11})$ whereas $o(A_{11})$ contains only one such (i.e. itself); this will contradict the lattice preserving nature of (+). Hence (ii) and the lemma.

Let A_{11} and A_{21} be elements of \underline{A}_{1} with $P_{1}(A_{11} \land A_{21}) = 0$. Since $\sigma(A_{11}) = \sigma(A_{21})$, it is obvious that the association mentioned in (ii) of lemma 9.2 preserves equivalence classes.

i.e. if for A_{11} we have A_{12} and A_{21} we have A_{22} through the association mentioned in (ii) of lemma 9.2, then $P_2(A_{12} \land A_{22}) = 0$. Thus we can unambiguously define a mapping T_1 from $\{[A_{11}]/A_{11} \in A_1, P_1(A_{11}) < \frac{1}{2}\}$ to $\{[A_{12}]/A_{12} \in A_2, P_2(A_{12}) < \frac{1}{2}\}$ using (ii) of lemma 9.2. T_1 is one to one is easy to check (using the fact that (4) preserves the metric). Working with (4) in the place of (4) we can conclude that T_1 is onto $\{[A_{12}]/A_{12} \in A_2, P_2(A_{12}) < \frac{1}{2}\}$. We note for future use that T_1 'preserves measure' in the sense if $T_1([A_{11}]) = [A_{12}]$ then $P_1(A_{11}) = P_2(A_{12})$. The next lemma is meant to show that T_1

Lemma 9.3 (i) Let $\mathbf{A}_{11} \in \underline{\mathbf{A}}_{1}$ be such that $P_{1}(\mathbf{A}_{11}) < \frac{1}{2}$. Let \mathbf{A}_{21} and \mathbf{A}_{31} be elements of $\underline{\mathbf{A}}_{1}$ satisfying $P_{1}(\mathbf{A}_{21}) = 0$, $P_{1}((\mathbf{A}_{21}) = 0)$ Let \mathbf{A}_{12} , \mathbf{A}_{22} and \mathbf{A}_{32} be elements of $\underline{\mathbf{A}}_{2}$ such that $T_{1}([\mathbf{A}_{j1}]) = [\mathbf{A}_{j2}]$, j = 1, 2, 3. Then $P_{2}(\mathbf{A}_{22}) = 0$ and

is well behaved.

$$P_2 ((A_{22} (_) A_{32}) \triangle A_{12}) = 0$$

(ii) Let $A_{O1} \in \underline{A}_{1}$ with $P_{1}(A_{O1}) < \frac{1}{2}$. Let

 $\{A_{n1}\}_{n \ge 1} \subseteq \underline{A}_{1}$ satisfying

 $P_1(A_{n1} - A_{01}) = 0$ and $P_1(A_{n1} - A_{n+1.1}) = 0$, $n \ge 1$.

Let A_{02} and $\{A_{n2}\}_{n\geq 1}$ be such that $T_1([A_{n1}]) = [A_{n2}], n \geq 0$.

Then, $P_2(A_{n2} - A_{02}) = 0$, $P_2(A_{n2} - A_{n+1,2}) = 0$ and

 $T_1([(n)A_{n1}]) = [(n)A_{n2}].$

<u>Proof</u>: (For (i)). A₁₁ & σ (A₂₁) V σ (A₃₁).

So, $A_{12} \in \sigma(A_{22}) \vee \sigma(A_{32})$ (recall that (+) preserves lattice operations and

 $T_1([A_{21}]) = [A_{22}] \Rightarrow (+)(\sigma(A_{21})) = \sigma(A_{22}) \text{ etc.}$

i.e. $A_{12} \in \sigma$ (the partition A_{22} A_{32} , A_{22} A_{32}^c , A_{22}^c A_{32}^c , A_{22}^c A_{32}^c).

 $P_{2}(A_{22}^{c}) = P_{1}(A_{11}) < \frac{1}{2}$. On the other hand,

$$P_{2}(A_{12} - (A_{22} A_{32} (_) A_{22} A_{32}^{c} (_) A_{22}^{c} A_{32})) = 0$$

$$i.e. P_{2}(A_{12} - (A_{22} (_) A_{32})) = 0$$

$$\Rightarrow P_{1}(A_{11}) = P_{2}(A_{12}) \leq P_{2}(A_{22} (_) A_{32})$$

$$= P_{2}(A_{22}) + P_{2}(A_{32}) - P_{2}(A_{22}(_) A_{32})$$

$$= P_{1}(A_{21}) + P_{1}(A_{31}) - P_{2}(A_{22}(_) A_{32})$$

$$= P_{1}(A_{21}) - P_{2}(A_{22}(_) A_{32}) \leq P_{1}(A_{11})$$

 \Rightarrow P₂ (A₂₂ () A₃₂) = 0. Also the fact P₂(A₁₂) = P₂(A₂₂ () A₃₂) (which follows from the above string of inequalities) ensures P₂ (A₁₂ \triangle (A₂₂ () A₃₂)) = 0. Hence (i).

Let now Λ_{11} , Λ_{21} , Λ_{12} and Λ_{22} be such that $P_1(\Lambda_{21}) < \frac{1}{2}$, $P_1(\Lambda_{11} - \Lambda_{21}) = 0$ and $T_1([\Lambda_{j1}]) = [\Lambda_{j2}]$, j=1,2. From (i) we have $P_2(\Lambda_{12} - \Lambda_{22}) = 0$. Now the proof of (ii) is trivial. Hence the lemma.

Loosely speaking Lemma 9.3 asserts that T_1 is 'monotone' and 'additive' (- part (i) -) and T_1 preserves'increasing limits' and hence 'countably additive' (- part (ii) -). That is, T_1 is an isomorphism from $\underline{A}_1^{(P_1)}$ ($\underline{A}_{11}^{(P_1)}$ onto $\underline{A}_2^{(P_2)}$ ($\underline{A}_{11}^{(P_1)}$) whenever $\underline{A}_1^{(P_1)}$ ($\underline{A}_{11}^{(P_1)}$) whenever $\underline{A}_1^{(P_1)}$ ($\underline{A}_1^{(P_1)}$) and $\underline{A}_2^{(P_2)}$ ($\underline{A}_1^{(P_1)}$) whenever $\underline{A}_1^{(P_1)}$ ($\underline{A}_1^{(P_1)}$) and $\underline{A}_2^{(P_2)}$ ($\underline{A}_1^{(P_1)}$) whenever $\underline{A}_1^{(P_1)}$ ($\underline{A}_1^{(P_1)}$) and $\underline{A}_2^{(P_2)}$ ($\underline{A}_1^{(P_1)}$) whenever $\underline{A}_1^{(P_1)}$ ($\underline{A}_1^{(P_1)}$) and $\underline{A}_1^{(P_1)}$ onto

Now we would like to dispense with two trivial cases. Let P_1 be completely atomic on $\underline{\mathbb{A}}_1$ with exactly one atom. Then $\underline{\mathbb{S}}_1$ contains only one element. So, $\underline{\mathbb{S}}_2$ will contain

only one element. This will imply that P2 is completely atomic on $\underline{A}_{\!\!\!2}$ with only one atom. Then the answer to our question is immediate. Next, let $\underline{\underline{P}}_1$ be completely atomic on $\underline{\mathbb{A}}_1$ with exactly two atoms (say \mathbb{A}_{11} , \mathbb{A}_{21}). Then $\underline{\mathbb{S}}_1$ has two elements. This will imply that $\c S_{lpha}$ has only two elements. So, P_2 will be completely atomic on $\underline{\mathbb{A}}_2$ with exactly two atoms (say A_{22} A_{22}). If $P_1(A_{11}) < \frac{1}{2}$, using T_1 it is seen that either $^{\Lambda}_{12}$ or $^{\Lambda}_{22}$ has $^{P}_{2}$ measure equal to $P_{\gamma}(A_{\gamma})$; then the affirmative answer to our question follows immediately. The case is same if $P_1(A_{21}) < \frac{1}{2}$. If $P_{1}(A_{11}) = P_{1}(A_{21}) = \frac{1}{2}$ then using the metric preserving nature of (+), it is evident that $P_2(\Lambda_{12}) = P_2(\Lambda_{22}) = \frac{1}{2}$. Again, the answer is immediate. So, hereafter we will exclude these two cases from our consideration.

We are set to define the mapping T from $A_1(P_1)$ onto $A_2(P_2)$. Get a partition $A_{11}, A_{21}, \ldots, A_{m1}$ of A_{m1} by A_{m1} sets as follows. A_{11} is the atom with the largest P_1 -measure, if there are no atoms take $A_{11} = \emptyset$. The remaining sets $A_{21}, A_{31}, \ldots, A_{m1}$ are so obtained that $P_1(A_{j1}) < \frac{1}{2}$, $j=2,3,\ldots,m$. This can be fone since we have excluded the case when P_1 is completely atomic with only one or two atoms. Let $A_{22}, A_{32}, \ldots, A_{m2}$ be A_{m2} sets such that $A_{j2} \in T_1(A_{j1}), j=2,3,\ldots,m$.

Now we claim that $P_2(A_{j2}(\vec{)}) = 0$, $2 \le j \ne j' \le m$. For if $D \in T_1^{-1}([A_{j2}(\overline{D}), A_{j+2}])$ then $P_1(D - A_{j1}) = 0$ and $P_1(D - A_{i+1}) = 0$. (This is a consequence of (i) of lemma 9.3) But $P_1(A_{j1}(\bar{D}) A_{j+1}) = 0$. So, $P_1(D) = 0$. Thus $P_2(A_{j2}(\overline{}) A_{j'2}) = 0$, $2 \le j \ne j' \le m$. Thus A_{j2} , A_{j3} , ..., A_{jm} are essentially disjoint. Let $\Lambda_{12} = (1)_2 - (1)_2 = (1)_2$ We presently show that $P_2(A_{12}) = P_1(A_{11})$ and if A_{11} is nonempty then A_{12} is also an atom. That $P_2(A_{12}) = P_1(A_{11})$ follows as $P_2(A_{j2}) = P_1(A_{j1})$, j = 2, 3, ..., m. Let if possible A_{j1} be nonempty and \mathbb{A}_{12} not an atom. Choose C (\mathbb{A}_{12} with $0 < P_2(C) < \min\{\frac{1}{2}, P_2(A_{12})\}$. Let $D \in T_1^{-1}[C]$. Since $P_1(A_{11}) = P_2(A_{12}) > P_2(C) = P_1(D)$ and A_{11} is an atom $P_1(D()) = 0$, So, $P_1(D()) = 0$ for some j = 2, 3, ..., m. Let: $E \in T_1[D(\overline{\ }) A_{j1}]$. Then $P_2(E - A_{j2}) = 0$ and $P_2(E - C) = 0$. (T[A₁] = [A₂] and T[D] = [C]). But $P_2(A_2$; () C) = 0. So, $P_2(E) = 0$ - a contradiction to the fact $P_1(D(T) A_{i1}) > 0$ and consequently $P_2(E) > 0$. Hence Λ_{12} is an atom.

Define T by

$$T([D]) = (\underbrace{)}_{j=1}^{m} T_{1}([D(\widehat{)} A_{j1}]), D \in \underline{A}_{1}$$

(Here interpret $T_1([A_{11}]) = [A_{12}]$. This interpretation is consistent if $P_1(A_{11}) < \frac{1}{2}$).

Lemma 9.3 and the observations we have made just above ensure that T is an isomorphism. In order not to lengthen an already long proof we omit the details, which are in any case easy to check.

Lemma 9.4 Let (+) be the map induced by T. Then for any D ϵ \underline{A}_1 ,

$$(+)$$
 $(\sigma_1D_1) = (+) (\sigma_1D_1).$

Proof: 1° . Let $P_1(D) < \frac{1}{2}$. Then

$$T([D]) = \bigoplus_{j=1}^{m} T_{1}([D (\overline{\ }) A_{j1}])$$

$$= T_{1}(\bigoplus_{j=1}^{m} [D (\overline{\ }) A_{j1}]) = T_{1}([D]).$$

Let $E \in T_1([D])$. By the definition of T_1 we have $(+)(\sigma\{D\}) = \sigma\{E\}$. On the other hand $E \in T([D])$ and by the definition of (+) we have (+) $(\sigma\{D\}) = \sigma\{E\}$.

 ∞ , $(+)(\sigma_{\{D\}}) = (+)(\sigma_{\{D\}})$.

 2° . Let E_1 , E_2 be two elements of \underline{A}_2 satisfying

(i) $P_2(E_1(\overline{}) E_2) = 0$, (ii) $0 < P_2(E_1) < \frac{1}{2}$ and

(iii) $P_2(E_1(\underline{}) E_2) = \frac{1}{2}$. Consider $\sigma\{E_1, E_2\}$. It is clear

that it contains only five elements of \S_2

viz $\sigma(\Sigma_2)$, $\sigma(E_1)$, $\sigma(E_2)$, $\sigma(E_1)$ and $\sigma(E_1, E_2)$.

So, the only element of \underline{S}_2 contained in $\sigma_1 E_1$, E_2 , and is at a distance $\frac{1}{2}$ from $\sigma_2 (\underline{\Gamma}_2)$ is $\sigma_3 (E_1 (\underline{\Gamma}_2)) E_2$.

 3° . Let $P_{\underline{1}}(D) = \frac{1}{2}$. Let B be an $\underline{A}_{\underline{1}}$ set with B \subset D and $0 < P_{\underline{1}}(B) < P_{\underline{1}}(D)$.

σ{D} (σ{B} Vσ{D - B} . So,

 $(+)(\sigma\{D\})$ (= $(+)(\sigma\{B\})$ V (+) $(\sigma\{D-B\})$ and

 $(+)(\alpha \{D\}) \subset (+)(\alpha \{B\}) \vee (+)(\alpha \{D - B\})$

 $= (+) (\alpha \{B\}) \quad \forall (+) (\alpha \{D - B\})$

(- note that $P_1(B) < \frac{1}{2}$, $P_1(D - B) < \frac{1}{2}$ and use 1° -).

So, $(+)(\sigma\{D\})$ and $(+)(\sigma\{D\})$ are elements of \underline{S}_2 contained in $(+)(\sigma\{B\})$ \forall $(+)(\sigma\{D-B\})$.

Let $E_1 \in T_1[B]$ and $E_2 \in T_1[D - B]$. Then

(i) $P_2(E_1(-)) = 0$, (ii) $0 < P_2(E_1) < \frac{1}{2}$, $0 < P_2(E_2) < \frac{1}{2}$ and (iii) $P_2(E_1(-)) = \frac{1}{2}$. Also, by the definition of T_1 , we have

$$(+) (\sigma \{B\}) \ V (+) (\sigma \{D - B\}) = \sigma \{E_1\} \ V \sigma \{E_2\}$$
.

Observe also that $d_1(\sigma_{\{(1)\}}, \sigma_{\{D\}}) = \frac{1}{2}$; the fact that (+) and (+) are isometries implies,

 $d_2(\sigma(\Omega_2), (+)(\sigma(D))) = \frac{1}{2}$ and $d_2(\sigma(\Omega_2), (+)(\sigma(D))) = \frac{1}{2}$.

Thus $(+)(\sigma\{D\})$ and $(+)(\sigma\{D\})$ are two elements of \underline{S}_2 , contained in $\sigma\{E_1\}$ $V \sigma\{E_2\}$ and each is at a distance $\frac{1}{2}$ from $\sigma\{C_2\}$. So, from 2° , $(+)(\sigma\{D\}) = (+)(\sigma\{D\})$.

 4° . Let $D \in \underline{A}_{1}$. If $P_{1}(D) \neq \frac{1}{2}$, either

 $P_1(D) < \frac{1}{2}$ or $P_1(D^c) < \frac{1}{2}$. Then an application of 1^o gives us $(+)(\sigma\{D\}) = (+)(\sigma\{D\})$. Now, let $P_1(D) = \frac{1}{2}$. Since we have excluded the case of exactly two atoms either D or D^c is not an atom; say D is not. Then D contains an Δ_1 set B with $P < P_1(B) < P_1(D)$. Now an application of 3^o gives us $(+)(\sigma\{D\}) = (+)(\sigma\{D\})$.

. Hence the lemma.

Lemma 9.5 For any $\underline{B} \in \underline{S}_1$, $(+)(\underline{B}) = (+)(\underline{B})$.

<u>Proof</u>: Observe that $\underline{B} = V \sigma \{D\}$. Recall that (+) and (+) are complete lattice isomorphisms. So,

$$(+)(\underline{B}) = (+)(\underbrace{V}_{D \in \underline{B}} \sigma \{D\}),$$

$$= \underbrace{V}_{D \in \underline{B}} (+)(\sigma \{D\})$$

$$= \underbrace{V}_{D \in \underline{B}} (+)(\sigma \{D\}) \quad (by lemma 9.4)$$

$$= (+)(\underbrace{V}_{D \in \underline{B}} \sigma \{D\}) = (+)(\underline{B}).$$

Hence the lemma.

Thus we have

Theorem 9.6 Every isometry (+) from $\subseteq (\cap_1, A_1, P_1)$ onto $\subseteq (\cap_2, A_2, P_2)$ preserving the lattice operations is induced by an isomorphism between $A_1(P_1)$ and $A_2(P_2)$.

CHAPTER 3

ON A GENERALIZATION, DUE TO BLAKE, OF MARTINGALES

1. Introduction

Let (\subseteq) , A P) be a probability space and $\{E_n\}_{n\geq 1}$ and increasing family of sub σ -algebras of A. Let $\{X_n\}_{n\geq 1}$ be a stochastic process adapted to $\{E_n\}_{n\geq 1}$; i.e. each X_n is E_n measurable. Following Blake [5], we refer to $\{X_n\}_{n\geq 1}$ as a game and define

Definition: The game $\{x_n\}_{n\geq 1}$ will be said to become fairer with time if for every $\epsilon > 0$,

 $P[|E(X/F_m) - X_m| > \epsilon] \to 0 \text{ as } n, m \to \infty \text{ with } n \ge m.$

(In the above definition the symbol $E(X_n/\underline{F}_m)$ stands for the conditional expectation of X_n given \underline{F}_m . We have used this notation for conditional expectation, rather than the one used in Chapter 1, to confirm with the usage in [5]). Any martingale is, trivially, a fairer with time game and thus this concept generalizes that of martingales.

Blake in [5] gives two convergence theorems for uniformly integrable fairer with time processes. Below we quote those results. Let $\{\alpha_n\}_{n\geq 1}$ be a monotonic sequence decreasing to zero with finite sum. The game $\{X_n\}_{n\geq 1}$ may be decomposed

with respect to $\{\alpha_n\}_{n \ge 1}$ as

$$\mathbf{x}_{\mathbf{n}} = \mathbf{Y}_{\mathbf{n}} - \mathbf{Z}_{\mathbf{n}} \tag{1.1}$$

where $\{Y_n\}_{n\geq 1}$ and $\{Z_n\}_{n\geq 1}$ are defined inductively by

$$Y_1 = X_1$$

$$Y_n = Y_{n-1} + [X_n - E(X_n^*/F_{n-1})] + \alpha_{n-1}$$

$$Z_n = Z_{n-1} + [X_{n-1} - E(X_n/F_{n-1})] + \alpha_{n-1}$$
(1.3)

The decomposition of $\{X_n\}_{n\geq 1}$ according to (1.1) - (1.3) will be called a Doob-like decomposition. It can be verified that $\{Y_n\}_{n\geq 1}$ of the decomposition, is a submartingale with respect to $\{F_n\}_{n\geq 1}$. Define the collection of sets $\{B_{n,m}\}_{n\geq 1}$ for $m=1,2,\ldots$ and $n\geq m$ by

 $B_{n,m}^{\alpha} = \{|E(X_n/F_m) - X_m| > \alpha_m\}$. Now we state the convergence theorems in [5].

Theorem 1. Let $\{X_n\}_{n\geq 1}$ be a uniformly integrable game and $\{Y_n\}_{n\geq 1}$, the submartingale associated with its Doob-like decomposition, be uniformly dominated in absolute value by an element of $L_1(\Omega)$, Δ , P). Suppose for every $\delta > 0$ there exists an integer N(δ) such that

Then, there exists a function X in $L_1(\Omega)$, A, P) such that $\lim_{n\to\infty}\int_{0}^{\infty}|X_n-X|\,dP=0$.

Theorem 2. Let $\{X_n\}_{n \ge 1}$ be a uniformly integrable game satisfying (1.4) and (1.5). Then, there exists some constant C such that $\lim_{n \to \infty} \int_{0}^{\infty} X_n dP = C$.

The condition (1.4) implies the game $\{X_n\}_n \ge 1$ is fairer with time. To see this fix $\epsilon > 0$. Let $\delta > 0$ be given. Choose m_0 so large that $\alpha_{m_0} < \epsilon$ (this can be done as $\alpha_{m_0} \downarrow 0$) and $m_0 \ge N(\delta)$. Now,

$$\{|E(X_n/\underline{F}_m) - X_m| > \varepsilon\} \subset \{|E(X_n/\underline{F}_m) - X_m| > \alpha_m\}$$
for all $n \ge m \ge m_0$.

So,
$$P[|E(X_n/\underline{F}_m) - X_m| > \epsilon] \le P[|E(X_n/\underline{F}_m) - X_m| > \alpha_m]$$

$$= P[B_{n,m}^{\alpha}] < \delta \quad \text{whenever}$$

$$n \ge m \ge m_0.$$

Since δ is arbitrary it follows $P[|E(X_n/\underline{F}_m) - X_m| > \epsilon] \rightarrow 0 \text{ as } n, m \rightarrow \infty \text{ with } n \geq m.$

Thus the above theorems give a set of sufficient conditions for a uniformly integrable fairer with time game to converge. In this chapter we show that these sufficient conditions are not needed; in fact, we show that any uniformly integrable, fairer with time game converges in L_1 .

2. Main theorems

Theorem 2.1 Any uniformly integrable fairer with time game $\{x_n\}_{n\geq 1}$ converges in L_1 .

<u>Proof</u>: To facilitate understanding, we break up the proof into a few important steps numbered (S1) through (S5). For every m and $n \ge m$ define $Y_{m,n} = E(X_n/F_m)$. Let Γ stand for the family $\{Y_{m,n}, \text{ for all } m \text{ and } n \ge m\}$.

(S1) Γ is uniformly integrable.

Since $\{X_n\}_n \ge 1$ is uniformly integrable there exists a function f defined on the non-negative real axis which is positive, increasing, and convex, such that

$$\lim_{t \to \infty} \frac{f(t)}{t} = +\infty$$

and $\sup_{n} \mathbb{E}[f \circ |X_{n}|] < \infty$. (see [16, II T 22]).

Now,
$$E[f \circ |Y_{m,n}|] = E[f \circ |E(X_n/\underline{F}_m)|]$$

$$\leq E[f \circ E(|X_n|/\underline{F}_m)] \text{ (since f is increasing)}$$

$$\leq E[E(f \circ |X_n|/\underline{F}_m)]$$

$$= E[f \circ |X_n|].$$

Therefore $\sup_{Y_{m,n} \in \Gamma} \mathbb{E}[f \circ |Y_{m,n}|] \leq \sup_{n} \mathbb{E}[f \circ |X_{n}|] < \infty$.

Another application of II T 22 of [16] ensures that \(\text{is} \) uniformly integrable. Hence (S1).

(S2) Given $\epsilon > 0$, there exists M such that for all $m \ge M$, one has $E(|X_m - Y_{m,n}|) \le 2\epsilon$ for all $n \ge m$.

Since Γ is uniformly integrable given $\epsilon > 0$ there exists $\delta > 0$ such that $P(A) < \delta$ implies

 $\int\limits_{A} |Y_{m,n}| \ dP \leq \frac{\epsilon}{2} , \text{ for all } Y_{m,n} \epsilon \quad \text{Choose M so large that}$ $m \geq M \text{ and } n \geq m \text{ implies } P[|X_m - E(X_n/\underline{F}_m)| > \epsilon] < \delta \text{ . Then,}$ it is not difficult to see that

 $E[|X_m - Y_{m,n}|] \le 2\varepsilon$ for all $m \ge M$ and $n \ge m$. Hence (S2).

(S3) For every fixed m, the sequence $\{Y_{m,n}\}$ converges in L_1 to an \underline{F}_m measurable random variable Z_m .

Let $m \leq n \leq n^{\dagger}$.

$$E[|Y_{m,n} - Y_{m,n}|] = E[|E(X_{n}/\underline{F}_{m}) - E(X_{n}/\underline{F}_{m})|]$$

$$= E[|E(X_{n} - X_{n},/\underline{F}_{m})|]$$

$$= E[|E(\{X_{n} - X_{n},/\underline{F}_{n}\}/\underline{F}_{m})|]$$

$$\leq E[E(\{|E(X_{n} - X_{n},/\underline{F}_{n})|\}/\underline{F}_{m})]$$

$$= E[|E(X_{n} - X_{n},/\underline{F}_{n})|]$$

$$= E[|X_{n} - Y_{n,n},|]$$

Now from (S2) it follows that given $\varepsilon > 0$ for all sufficiently large n and n'

$$E[|Y_{m,n} - Y_{m,n}|] \le E[|(X_n - Y_{n,n}|)] \le 2\varepsilon$$
.

Hence, for m fixed, the sequence $\{Y_{m,n}\}$ is Cauchy in the L_1 -norm. So, there exists, an integrable random variable Z_m , such that, $Y_{m,n} \xrightarrow{L_1} Z_m$. Without loss in generality we can take Z_m to be \underline{F}_m measurable. (Note that each $Y_{m,n}$ is \underline{F}_m measurable and there is a subsequence $\{Y_{m,n}\}$ converging almost surely to Z_m). Hence (S3).

(S4) $\{Z_m, \underline{F}_m\}_{m \ge 1}$ is an uniformly integrable martingale.

The fact that $\{Z_m\}_m \ge 1$ is uniformly integrable follows trivially because the closure in L_1 of a uniformly integrable collection is uniformly integrable. (See [16, II T 20]). To show $\{Z_m, \underline{F}_m\}$ is a martingale it is enough to show that for every m,

$$\begin{split} & E(Z_{m+1}/\underline{F}_{m}) = Z_{m} \quad \text{a.s. Since} \\ & E[|E(Y_{m+1,n}/\underline{F}_{m}) - E(Z_{m+1}/\underline{F}_{m})|] \\ & = E[|E((Y_{m+1,n} - Z_{m+1})/\underline{F}_{m})|] \\ & \leq E[E[|Y_{m+1,n} - Z_{m+1})/\underline{F}_{m}]] \\ & = E[|Y_{m+1,n} - Z_{m+1}|] \rightarrow 0 \quad \text{as} \quad n \rightarrow \infty \,, \end{split}$$

there exists a subsequence n! of $\{n: n \ge m\}$ such that

$$\mathbb{E}(\mathbb{Y}_{m+1,n},/\underline{\mathbb{F}}_{m}) \stackrel{a.s.}{\rightarrow} \mathbb{E}(\mathbb{Z}_{m+1}/\underline{\mathbb{F}}_{m}).$$

We can assume (- if necessary, by choosing a further subsequence,-) that $Y_{m,n}$: $\stackrel{a.s.}{\longrightarrow} Z_m$.

Now,
$$E(Z_{m+1}/\underline{F}_m) = \lim_{n \to \infty} E(Y_{m+1,n}, /\underline{F}_m)$$
 a.s.

$$= \lim_{n \to \infty} E(\{E(X_n, /\underline{F}_{m+1})\}/\underline{F}_m) \text{ a.s.}$$

$$= \lim_{n \to \infty} E(X_n, /\underline{F}_m) \text{ a.s.}$$

$$= \lim_{n \to \infty} Y_{m,n}, \text{ a.s.}$$

$$= Z_m \text{ a.s.}$$

Hence (S4).

(S5) $\{X_n\}_{n \geq 1}$ converges in L_1 .

Since $\{Z_n, E_n\}_{n \ge 1}$ is an uniformly integrable martingale, there exists an integrable random variable Z_n such that

 $Z_n \xrightarrow[n \to \infty]{L_1} Z_\infty$. We shall show that $X_n \xrightarrow[n \to \infty]{L_1} Z_\infty$. From (S3) and (S2) it is easy to check that given $\varepsilon > 0$ there exists M such that for all $m \ge M$

 $|X_m - Z_m| dP \le 2\varepsilon$. Therefore, for sufficiently large m, $|X_m - Z_m| dP \le \int |X_m - Z_m| dP + \int |Z_m - Z_\infty| dP \le 3\varepsilon$, say. Hence (S5) and the theorem.

Since any game (stochastic process) $\{x_n\}_n \ge 1$ converging in L_1 can be taken to be a game fairer with time, by setting $F_n = A$ for all n, we get the following corollary.

Corollary 2.1 Let $\{X_n\}_{n \ge 1}$ be a game. It converges in L_1 if and only if it is uniformly integrable and fairer with time with respect to some increasing family of sub σ -algebras

 $\{F_n\}_{n \geq 1}$ to which it is adapted.

Let p > 1.

Theorem 2.2 Let $\{X_n\}_{n \ge 1}$ be a fairer with time game with $\{|X_n|^p\}_{n \ge 1}$ uniformly integrable. Then $\{X_n\}_{n \ge 1}$ converges in L_p .

<u>Proof</u>: Noting that the function f defined on the non-negative real axis by $f(t) = t^p$ is positive, increasing and convex and $\lim_{t\to\infty} \frac{f(t)}{t} = +\infty$, in view of II T 22 of [16], it is clear

that $\{X_n\}_n \ge 1$ is uniformly integrable. Hence by theorem 2.1 it converges in L_1 ; in particular, $\{X_n\}_n \ge 1$ converges in probability. Therefore $\{X_n\}_n \ge 1$ converges in L_p . (See proposition II 6.1 of [20]).

Corollary 2.2 The game $\{X_n\}_{n\geq 1}$ converges in L_p if and only if $\{|X_n|^p\}_{n\geq 1}$ is uniformly integrable and $\{X_n\}_{n\geq 1}$ is fairer with time with respect to some increasing family of sub σ -algebras $\{F_n\}_{n\geq 1}$ to which it is adapted.

CHAPTER 4

ON A CONJECTURE ON SINGULAR MARTINGALES

1. Introduction

$$\mu_{\mathbf{n}}(\mathbf{A}) = \int_{\mathbf{A}} \mathbf{X}_{\mathbf{n}} d\mathbf{P}, \quad \mathbf{A} \in \underline{\mathbf{F}}_{\mathbf{n}}.$$
 (1.1)

In [14], Luis Baez-Duarte introduced the following definitions.

Definition 1.2 The martingale $\{X_n, \underline{F}_n\}_n \geq 1$ is measure—dominated if there is a finite measure μ on \underline{B} whose restriction to each \underline{F}_n coincides with the measure μ_n defined in (1.1). In such a case μ is said to dominate the martingale.

In view of our assumption $\underline{F}_{\infty} = \underline{A}$, if there is a measure μ dominating the martingale $\{X_n, \underline{F}_n\}_n \ge 1$ then it is unique.

Definition 1.3 Let $\{X_n, E_n\}_{n \geq 1}$ be a measure-dominated martingale. It is said to be a singular martingale if the dominating measure μ is singular with respect to P.

Let for each n, S_n stand for the extended real line $[-\infty,\infty]$ and \underline{B}_n for the Borel σ -algebra of S_n . Denote by S the cartesian product $\prod_{n=1}^{\infty} S_n$ and by Y_n the projection from S to (S_n,\underline{B}_n) . Let the σ -algebras generated by $\{Y_1,Y_2,\ldots,Y_n\}$ and $\{Y_n,n\geq 1\}$ be denoted respectively by \underline{B}^n and \underline{B} . In the remainder of this chapter, the symbols $S,S_n,\underline{B}^n,\underline{B}$ and Y_n have the same meaning as given in this paragraph.

Given any sequence $\{X_n\}_{n\geq 1}$ of integrable random variables on $(\ \underline{\ \ \ })$, \underline{A} , \underline{P}) by canonical mapping we mean the measurable mapping \underline{T} from $(\ \underline{\ \ \ })$, \underline{A}) to $(\underline{S}$, \underline{B}) given by

$$T(\omega) = (X_1(\omega), X_2(\omega), \dots, X_n(\omega), \dots), \forall \omega \in \Omega$$
.

It is easy to see that if $\{X_n, \underline{F}_n\}_n \geq 1$ is a martingale on (\underline{C}) , \underline{A} , \underline{P}) then $\{Y_n, \underline{B}^n\}_n \geq 1$ is a martingale on (s, \underline{B} , \underline{P}) (- where \underline{T} is, of course, the canonical mapping -) with the property that $\{X_n\}_n \geq 1$ and $\{Y_n\}_n \geq 1$ are equivalent. We call $\{Y_n, \underline{B}^n\}_n \geq 1$ the canonical martingale associated with $\{X_n, \underline{F}_n\}_n \geq 1$. Note that for

various $\{X_n\}$'s what varies is the measure P o T⁻¹ and not $\{Y_n\}_{n \ge 1}$

It is not difficult to prove (see, for instance, theorem 3 l of [14]) that, if $\{X_n, F_n\}_{n \geq 1}$ is a measure-dominated martingale then the canonical martingale associated with it is also measure-dominated. In [14], Baez-Duarte conjectured that if $\{X_n, F_n\}_{n \geq 1}$ is a singular martingale on $(\Gamma), A$, P) then the corresponding dominating measure for the canonical martingale is concentrated on the boundary of Γ We disprove this with the help of an example.

2. An example

Before proceeding with the construction of the example we would like to observe the following.

Let \subseteq S, $\underline{A} = \underline{B}$ and $X_n = Y_n$ for all n. Then the canonical mapping from (\subseteq A) to (S, B) is just the identity mapping. So $\underline{T} \subseteq \underline{S}$ and hence the boundary of $\underline{T} \subseteq \underline{S}$ is empty. Therefore if we could find a measure \underline{Q} on (S, B) such that \underline{Y}_n , \underline{B}_n becomes a singular martingale on (S, B, Q) we would have disproved the conjecture.

We achieve this in two stages.

2.1 Let $\subseteq \{1, 2, 3, ..., n, ..., \infty\}$. Let $= \{\emptyset, \subseteq \Sigma\}$ and, for each $n \ge 2$, $= \{0, 1\}$ be the \emptyset -algebra generated by the

partition $\{1\}$, $\{2\}$, ..., $\{n-1\}$, $\{n, n+1, \ldots, \infty\}$. The symbol $[n, \infty]$ will denote $\{n, n+1, \ldots, \infty\}$. As in section 1, the σ -algebra generated by $\binom{n}{n} = n$ is denoted by $\frac{\Delta}{n}$. Clearly, is just the class of all subsets of $\binom{n}{n}$. Let P be the probability measure defined on $\frac{\Delta}{n}$ by $P(\{n\}) = \frac{1}{2^n}$, $n \ge 1$ and $P(\{\infty\}) = 0$. Define a sequence $\binom{X}{n} = 1$ of integrable random variables on $\binom{n}{n} = 1$ by setting,

$$X_n(\omega) = \frac{1}{P[n, \infty]}$$
 if $\omega \ge n$

$$= 0 if \omega < n.$$

2.2 Consider the canonical martingale associated with the singular martingale constructed in 2.1. We show that it (the canonical martingale) is singular.

The following two facts are easily observed.

(i) Since the set (), defined in 2.1 is countable, T(A) is B measurable for every subset A of ().

(ii) Since
$$T(\omega) = (X_1(\omega), X_2(\omega), \dots, X_n(\omega), \dots)$$

= $(1, \frac{1}{P[2, \infty]}, \frac{1}{P[3, \infty]}, \dots, \frac{1}{P[\omega, \infty]}, 0, 0, \dots)$

I is one to one.

Now, denoting the measures, corresponding to the canonical martingale, defined in (1.1) by $\{\mu_n^i\}_{n\geq 1}$, we have

$$\mu_n^{\dagger}(\mathbf{A}) = \int_{\mathbf{A}} \mathbf{Y}_n \, d \, \mathbf{P} \, o \, \mathbf{T}^{-1} = \int_{\mathbf{T}^{-1}(\mathbf{A})} \mathbf{X}_n \, d\mathbf{P} = \mu_n(\mathbf{T}^{-1}(\mathbf{A})), \, \mathbf{Y} \, \mathbf{A} \, \epsilon \, \mathbf{B}^n.$$

That is $\mu_n^{!} = \mu_n$ o T^{-1} . Let $\mu^{!}$ be the probability measure concentrated at the point $T(\infty)$.

For any
$$A \in \underline{B}^n$$
, $\mu^!(A) = 1 \iff T(\infty) \in A$
$$\iff \mu_n(T^{-1}(A)) = 1$$

$$\iff \mu_n^!(A) = 1$$

So, for every n, the restriction of μ ' to \underline{B}^n coincides with μ_n^i . That is, μ^i is the dominating measure for the canonical martingale. Since, PoT⁻¹(T($\underline{\Gamma}$) - $\underline{\rho}$)) = P($\underline{\Gamma}$) - $\underline{\rho}$)= 1 and μ^i (T(∞)) = 1, the measures PoT⁻¹ and μ^i are singular.

Thus $\{Y_n, \underline{B}^n\}_{n \geq 1}$ is a singular martingale on (S, B, P o T⁻¹). Therefore, in view of the observation made at the beginning of this section, the conjecture made by Baez-Duarte is false.

3. Remark

A slight modification of the above example will show that even if boundary of T () is non-empty the conjecture is not true. We indicate below the modification to be effected without giving detailed proofs.

Let \bigcap and P be as in 2.1 and s ϵ S - T \bigcap .

Define S' = S - $\{s\}$ and take (\underline{B}^n) ' and \underline{B} ' respectively to be the trace σ -algebras of \underline{B}^n and \underline{B} on S'. Let, for each n, Y_n^i be the restriction of Y_n to S'. $\{Y_n^i, (\underline{B}^n)^i\}_n \ge 1$ is a singular martingale on $(S^i, \underline{B}^i, P \circ T^{-1})$ with μ^i (defined in 2.2) as the dominating measure. Now, since $T(S^i) = S - \{s\}$, the boundary of $T(S^i) = s$. But, the dominating measure corresponding to the canonical martingale associated with $\{Y_n^i, (\underline{B}^n)^i\}_n \ge 1$ is concentrated at the point $\{1, \frac{1}{P[2,\infty]}, \cdots, \frac{1}{P[n,\infty]}, \cdots\}$; and this point is obviously different from s.

REFERENCES

- [1] K.P.S. Bhaskara Rao and M. Bhaskara Rao A remark on nonatomic measures, The Annals of Mathematical Statistics, 43, 1972, pp. 369-370; Correction, p. 1380.
- [2] K.P.S. Bhaskara Rao and M. Bhaskara Rao Δ (μ), Indian Statistical Institute, Tech. Rep. no. Math-Stat./
- [3] M. Bhaskara Rao Contributions to measure theory, Thesis submitted to Indian Statistical Institute, 1971.
- [4] D. Blackwell On a class of probability spaces, Proceedings of the Third Berkeley Symposium on Mathematical Statistics and Probability, University of California Press, 1956, Vol. II, pp. 1-6.
- [5] L. H. Blake A generalization of martingales and two consequent convergence theorems, Pacific J. Math. 35, 1970, pp. 279-283.
- [6] E.D. Bolker Functions resembling quotients of measures, Trans. Amer. Math. Soc., 124, 1966, pp. 292-312.
- [7] E.S. Boylan Equiconvergence of martingales, The Annals of Mathematical Statistics, 42, 1972, pp. 552-559.

- [8] P.R. Halmos The range of a vector measure, Bull. Amer. Math. Soc. 54, 1948, 416-421.
- [9] P.R. Halmos Measure Theory, D. Van Nostrand Company Inc.,
- [10] F. Hausdorff Set Theory, Chelsea Publishing Company, New York, 1957.
- [11] K. Kuratowski Topology, Vol. I, Academic Press, 1966.
- [12] K. Kuratowski Topology, Vol. II, Academic Press, 1968.
- [13] J. Lindenstrauss A short proof of Liapounoff's convexity theorem, J. Math. Mech. 15, 1966, pp. 971-972.
- [14] Luis Baez-Duarte Another look at martingale theorem,

 Journal of Math. Analysis and Applications, 23,

 1968, pp. 551-557.
- [15] E. Marczewski and R. Sikorski On isomorphism types of measure algebras, Fund, Math., 38, 1951, pp.92-98.
- [16] P.A. Meyer Probability and Potentials, Blaisdell
 Publishing Company, Waltham, Massachusetess, 1966.
- [17] E. Michael Topologies on spaces of subsets, Trans. Amer. Math. Soc., 71, 1951, pp. 152-182.
- [18] J.I. Nagata Modern dimension theory, North Holland Publishing Company, 1965.

- [19] J. Neveu Atomes conditionnels d'espaces de probabilite et theorie de l'information, Symposium on Probability Methods in Analysis, Springer-Verlag lecture notes in mathematics, 31, 1967, pp. 256-271.
- [19a] J. Neveu Note on the tightness of the metric on the set of complete sub o-algebras of a probability space,

 The Annals of Mathematical Statistics, 43, 1973.
- [20] J. Neveu Mathematical Foundations of the Calculus of Probability, Holden-Day, Inc., 1965.
- [21] D. Ramachandran Mixtures of perfect probability measures,

 The Annals of Probability (to appear).
- [22] V.A. Rohlin On the fundamental ideas of measure theory,

 Mat. Sb., 25 (67) (1949), 107-150 (in Russian),

 Amer. Math. Soc. Transl. Series 1, 10, 1-54.
- [23] R. Subramanian On a generalization of martingales due to Blake, Pacific J. Math. (1973).
- [24] R. Subramanian On a conjecture on singular martingales, J. Math. Anal. Appl. (1973).