# Certain Functional of Probability Measures on Hilbert Spaces 

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## §1. Introduction and results

Let $E$ be a real separable Hilbert space with inner product <,> and $\mathscr{E}$ the $\sigma$-algebra of all Borel subsets of $E$. We denote by $\mathscr{P}$ the set of all probability measures $\mu$ on $(E, \mathscr{E})$ with a finite second moment; $\int\|x\|^{2} d \mu(x)<\infty$. For each $\mu \in \mathscr{P}$ there exist a vector $m$ (mean vector) and a bounded linear operator $V$ (covariance operator) with $\left.\int\langle x, u\rangle d \mu(x)=<m, u\right\rangle$ and $\int\langle x-m, u\rangle\langle x-m$, $v>d \mu(x)=\langle V u, v\rangle$ for all $u, v \in E$. Since the covariance operator is symmetric, non-negative and nuclear, we can find a unique Gaussian measure $\gamma_{\mu}$ on ( $E, \mathscr{E}$ ) which has the same mean vector and covariance operator as those of $\mu[4 ;$ p. 14 and p.18]. Let $\mathscr{M}(\mu)$ be the set of all probability measures $M$ on $(E \times E, \mathscr{E} \otimes \mathscr{E})$ with $M(A \times E)=\mu(A)$ and $M(E \times A)=\gamma_{\mu}(A)$ for all $A \in \mathscr{E}$. We consider a function: $M \rightarrow e[\mu ; M]=\iint\|x-y\|^{2} d M(x, y)$ on $\mathscr{M}(\mu)$, and define a functional $e$ on $\mathscr{P}$ by

$$
e[\mu]=\inf _{M \in \mathbb{N}(\mu)} e[\mu ; M] .
$$

The functional $e$ was first introduced by H. Tanaka in the case where $E$ is the onedimensional space and its basic properties were studied also by himself [5]. H. Murata and H. Tanaka [2] extended the results to the case of multi-dimensional Euclidean spaces.

The purpose of this paper is to show that some of their results can be extended to the case of Hilbert spaces, by the method similar to that of [2] with a slight simplification. That is, we shall prove:

Theorem 1. For each $\mu \in \mathscr{P}$ there exists an $M \in \mathscr{M}(\mu)$ with $e[\mu]=e[\mu$; $M]$ and such a measure $M$ has the form; $M(A \times B)=\gamma_{\mu}\left(f^{-1}(A) \cap B\right)$ for all $A$, $B \in \mathscr{E}$ with a Borel measurable mapping from $E$ into itself. Consequently $e[\mu]=\int\|f(y)-y\|^{2} d \gamma_{\mu}(y)$.

Theorem 2. Let $\mu_{1}$ and $\mu_{2}$ be measures in $\mathscr{P}$ and $\mu_{1} * \mu_{2}$ their convolution. Then

$$
e\left[\mu_{1} * \mu_{2}\right] \leqq e\left[\mu_{1}\right]+e\left[\mu_{2}\right],
$$

and the equality holds if and only if both $\mu_{1}$ and $\mu_{2}$ are Gaussian.
Using the results we shall prove also that a sequence of probability distributions of certain stochastic processes $X_{n}=\left(X_{n}(t)\right)_{0 \leqq t<1}$ converges to a Gaussian measure in $L_{2}[0,1)$.

## § 2. Lemmas

In this section we denote by $\mathscr{C}\left(E^{n}\right)$ the Banach space of all real valued, bounded and continuous functions on $E^{n}$ with the supremum norm; $\|\varphi\|_{\infty}=$ $\sup |\varphi(x)|$, and by $\mathscr{C}^{*}\left(E^{n}\right)$ the topological dual of $\mathscr{C}\left(E^{n}\right)$. Since, for each $M$ $\in \mathscr{M}(\mu)$, the function: $\varphi \rightarrow M(\varphi)=\int_{E^{2}} \varphi d M$ on $\mathscr{C}\left(E^{2}\right)$ is continuous and linear, we consider $\mathscr{M}(\mu)$ as a subset of $\mathscr{C}^{*}\left(E^{2}\right)$.

Lemma 1. For each $\mu \in \mathscr{P}$ there exists an $M \in \mathscr{M}(\mu)$ with $e[\mu]=e[\mu ; M]$.
Proof. We shall prove first that $\mathscr{M}(\mu)$ is a weakly compact subset of $\mathscr{C}^{*}\left(E^{2}\right)$. Let $U^{0}$ be the closed unit ball in $\mathscr{C}^{*}\left(E^{2}\right)$, which is known to be weakly compact. Since $\mathscr{M}(\mu)$ is contained in $U^{0}$, it is enough to show that $\mathscr{M}(\mu)$ is weakly closed. Let $M_{0}$ be an element in the weak closure of $\mathscr{M}(\mu)$. Then there is a net $\left(M_{\lambda}\right)_{\lambda \in \Lambda}$ in $\mathscr{M}(\mu)$ which converges weakly to $M_{0}$. It is easily seen that $M_{0}$ is linear and positive, and satisfies $M_{0}(1)=1$. For a given $\varepsilon>0$, we can find a compact subset $F$ of $E$ with $\mu(F) \geqq 1-\varepsilon / 2$ and $\gamma_{\mu}(F) \geqq 1-\varepsilon / 2$. Let $K=F \times F$. Then $M_{\lambda}(K) \geqq \mu(F)$ $-\gamma_{\mu}\left(F^{c}\right) \geqq 1-\varepsilon$ for all $\lambda \in \Lambda$. Therefore if a function $\varphi$ in $\mathscr{C}\left(E^{2}\right)$ vanishes on $K$, $\left|M_{0}(\varphi)\right|=\lim _{\lambda}\left|M_{\lambda}(\varphi)\right| \leqq \varepsilon\|\varphi\|_{\infty}$, which implies that $M_{0}$ is a Baire (hence Borel) probability measure on $E^{2}$. For any $\varphi \in \mathscr{C}(E)$, since $\varphi \circ \pi_{i} \in \mathscr{C}\left(E^{2}\right)(i=1,2)^{1)}$, we have $M_{0}\left(\varphi \circ \pi_{1}\right)=\lim _{\lambda} M_{\lambda}\left(\varphi \circ \pi_{1}\right)=\mu(\varphi)$ and similarly, $M_{0}\left(\varphi \circ \pi_{2}\right)=\gamma_{\mu}(\varphi)$, which shows that $M_{0}$ belongs to $\mathscr{M}(\mu)$. Thus $\mathscr{M}(\mu)$ is weakly closed. Now let $\varphi_{n}(x, y)$ $=\inf \left(n,\|x-y\|^{2}\right)$ and $\Phi_{n}(M)=\int_{E^{2}} \varphi_{n} d M$ for each $M \in \mathscr{M}(\mu)$. Then $\Phi_{n}$ are continuous on $\mathscr{M}(\mu)$ and $\Phi_{n} \uparrow e[\mu ; \cdot]$ as $n \rightarrow \infty$. Therefore $e[\mu ; \cdot]$ is lower semicontinuous on $\mathscr{M}(\mu)$ and hence, there is an $M \in \mathscr{M}(\mu)$ with $e[\mu]=e[\mu ; M]$.

From now on we use $\mathscr{M}_{0}(\mu)$ to denote the set of $M$, in $\mathscr{M}(\mu)$ and with $e[\mu]$ $=e[\mu ; M]$. We set $\Gamma=\left\{\left(x, y, x^{\prime}, y^{\prime}\right) \in E^{4}:\left\langle x-x^{\prime}, y-y^{\prime}\right\rangle \geqq 0\right\}$ and $\mathscr{M}_{\Gamma}(\mu)$ $=\{M \in \mathscr{M}(\mu): M \otimes M(\Gamma)=1\}$, where " $\otimes$ " denotes the direct product of measures.

Lemma 2. $\mathscr{M}_{0}(\mu) \subset \mathscr{M}_{\Gamma}(\mu)$ for all $\mu \in \mathscr{P}$.

[^0]Proof. To avoid the trivial case we assume $\int\|x-m\|^{2} d \mu(x)>0$, where $m$ is the mean vector of $\mu$. Under the assumption any $M \in \mathscr{M}(\mu)$ has no atomic point. Suppose that there is an $M \in \mathscr{M}_{0}(\mu)$ with $M \otimes M(\Gamma)<1$. There then exist $\mathscr{E}^{2}$ measurable sets $A$ and $B$ with $A \times B \subset \Gamma^{c}$ and $M(A) M(B)=M \otimes M(A \times B)>0$. We choose two compact sets $K$ and $K^{\prime}$ with $K \subset A$ and $K^{\prime} \subset B$ and $M(K) M\left(K^{\prime}\right)$ $>0$. The function: $\left(x, y, x^{\prime}, y^{\prime}\right) \rightarrow\left\langle x-x^{\prime}, y-y^{\prime}\right\rangle$ is continuous and strictly negative on $K \times K^{\prime}$ and hence, there is a constant $\delta>0$ with $\left\langle x-x^{\prime}, y-y^{\prime}\right\rangle$ $\leqq-\delta$ for all $\left(x, y, x^{\prime}, y^{\prime}\right) \in K \times K^{\prime}$. This implies also that $K$ and $K^{\prime}$ are disjoint. Now we consider an ergodic automorphism $T$ acting on ( $E^{2}, \mathscr{E}^{2}, M$ ) (such an automorphism always exists, for the space $\left(E^{2}, \mathscr{E}^{2}, M\right)$ is isomorphic $\bmod 0$ to the interval $[0,1)$ with the Lebesgue measure). By ergodicity there is an integer $n \geqq 0$ with $M\left(K \cap T^{-n} K^{\prime}\right)>0$. We set $C=K \cap T^{-n} K^{\prime}, C^{\prime}=T^{n} C$ and $D$ $=\left(C \cup C^{\prime}\right)^{c}$. The sets $C$ and $C^{\prime}$ are disjoint since they are subsets of $K$ and $K^{\prime}$ respectively. We define a mapping $S$ from $E^{2}$ into itself by

$$
S(z)=\left\{\begin{array}{lll}
\left(\pi_{1} T^{n} z, \pi_{2} z\right) & \text { for } & z \in C, \\
\left(\pi_{1} T^{-n} z, \pi_{2} z\right) & \text { for } & z \in C^{\prime}, \\
\left(\pi_{1} z, \pi_{2} z\right) & \text { for } & z \in D .
\end{array}\right.
$$

Let $M_{S}(A)=M\left(S^{-1}(A)\right)$ for all $A \in \mathscr{E}^{2}$. Then it is easy to see that $M_{S}$ belongs to $\mathscr{M}(\mu)$ and that

$$
\begin{aligned}
& e[\mu ; M]-e\left[\mu ; M_{S}\right] \\
&=\int_{E^{2}}\left\|\pi_{1}-\pi_{2}\right\|^{2} d M-\int_{E^{2}}\left\|\pi_{1}-\pi_{2}\right\|^{2} d M_{S} \\
&=-2 \int_{C}<\pi_{1} z-\pi_{1} T^{n} z, \pi_{2} z-\pi_{2} T^{n} z>d M(z) \\
& \geqq \delta M(C)>0,
\end{aligned}
$$

because $\left(z, T^{n} z\right) \in C \times C^{\prime} \subset K \times K^{\prime}$ whenever $z \in C$. This contradicts the assumption that $M$ belongs to $\mathscr{M}_{0}(\mu)$.

Let $M \in \mathscr{M}(\mu)$. A Markov kernel $P_{M}$, defined on $\mathscr{E} \times E$, is called a regular conditional probability of $M$ (with respect to $\gamma_{\mu}$ ) if $M(A \times B)=\int_{B} P_{M}(A, y) d \gamma_{\mu}(y)$ for all $A, B \in \mathscr{E}$.

Lemma 3. Let $M \in \mathscr{M}_{\Gamma}(\mu)$. Then there exist a set $N$ in $\mathscr{E}$ and a regular conditional probability $P_{M}$ of $M$ such that $\gamma_{\mu}(N)=0$ and, for any $y, y^{\prime} \in N^{c}$

$$
P_{M}(, y) \otimes P_{M}\left(, y^{\prime}\right)\left(\Gamma\left(y, y^{\prime}\right)\right)=1,
$$

where $\Gamma\left(y, y^{\prime}\right)=\left\{\left(x, x^{\prime}\right):\left(x, y, x^{\prime}, y^{\prime}\right) \in \Gamma\right\}$.
Proof. Let $\left(G_{n}\right)_{n \geqq 1}$ be a countable open base of $E$. For each $n \geqq 1$ we define a partition $\mathscr{B}_{n}=\left\{B_{n j}: 0 \leqq j \leqq 2^{n}-1\right\}$ of $E$ as follows; if $j=k_{1} 2^{n-1}+$ $k_{2} 2^{n-2}+\cdots+k_{n}$ is the binary expansion of $j, B_{n j}=G_{1}^{k_{1}} \cap G_{2}^{k_{2}} \cap \cdots \cap G_{n}^{k_{n}}$, where $G^{1}$ and $G^{0}$ are understood to be $G$ and $G^{c}$ respectively. We denote by $\mathscr{E}_{n}$ the $\sigma$-algebra generated by $\mathscr{B}_{n}$. Then the algebra $\mathscr{B}=U_{n \geqq 1} \mathscr{E}_{n}$ contains at most countable number of elements, and generates the $\sigma$-algebra $\mathscr{E}$. For a fixed $A$ $\in \mathscr{E}$ we define functions $Q_{M}^{n}(A$,$) on E$ by

$$
Q_{M}^{n}(A, y)=\left\{\begin{array}{lllll}
M\left(A \times B_{n j}\right) / \gamma_{\mu}\left(B_{n j}\right) & \text { if } & y \in B_{n j} & \text { and } & \gamma_{\mu}\left(B_{n j}\right)>0 \\
0 & \text { if } & y \in B_{n j} & \text { and } & \gamma_{\mu}\left(B_{n j}\right)=0
\end{array}\right.
$$

Then $\left(Q_{M}^{n}(A,), \mathscr{E}_{n}\right)_{n \geqq 1}$ is a martingale on the probability space $\left(E, \mathscr{E}, \gamma_{\mu}\right)$ and hence, there is a set $N_{A}$ with $\gamma_{\mu}\left(N_{A}\right)=0$ such that, for each $y \notin N_{A}, Q_{M}(A, y)=$ $\lim Q_{M}^{n}(A, y)$ exists. Now let $N=U_{A \in \mathscr{B}} N_{A}$. Then $\gamma_{\mu}(N)=0$ and, for each $y \notin N$, $Q_{M}(, y)$ is a finitely additive probability measure on $\mathscr{B}$. Using the injection: $x \rightarrow\left(I_{G_{n}}(x)\right)_{n \geqq 1}$ of $E$ into $\{0,1\}^{N_{0}}$, we can extend each $Q_{M}(, y)$ to a probability measure on $\mathscr{E}$, where $I_{G}$ denotes the indicator of a set $G$ (for the details, see [3]). We define probability measures $P_{M}(, y)$ by; $P_{M}(, y)=Q_{M}(, y)$ for $y \notin N$ and $P_{M}(, y)=\mu_{0}$ for $y \in N$, where $\mu_{0}$ is an arbitrary probability measure on $E$. Then $P_{M}$ is a regular conditional probability of $M$. We remark here that, for each $y \notin N$, the sequence of measures $\left(Q_{M}^{n}(, y)\right)_{n \geqq 1}$ converges weakly to $P_{M}(, y)$. In fact, for an $\varepsilon>0$ and an open set $G$ in $E$, we can find $A \in \mathscr{B}$ with $A \subset G$ and $P_{M}(G, y)$ $-\varepsilon \leqq P_{M}(A, y)$, for $\mathscr{B}$ contains the open base $\left(G_{n}\right)_{n \geqq 1}$. Hence $P_{M}(G, y)-\varepsilon$ $\leqq \lim Q_{M}^{n}(A, y) \leqq \liminf Q_{M}^{n}(G, y)$, and so, $P_{M}(G, y) \leqq \liminf Q_{M}^{n}(G, y)$, which implies that $\left(Q_{M}^{n}(, y)\right)$ converges weakly to $P_{M}(, y)$. By the assumption $Q_{M}^{n}($, $y) \otimes Q_{M}^{n}\left(, y^{\prime}\right)\left(\Gamma\left(y, y^{\prime}\right)\right)=1$ for almost all $\left(y, y^{\prime}\right)$ with respect to $\gamma_{\mu} \otimes \gamma_{\mu}$, however, for any $A \in \mathscr{E}$, since the function $Q_{M}^{n}(A$,$) is equal to a constant on each B_{n j}$, the equality holds for all $\left(y, y^{\prime}\right) \in E^{2}$. The sequence $\left(Q_{M}^{n}(, y)\right)_{n \geqq 1}$ converges weakly to $P_{M}(, y)$ for any $y \notin N$ and the set $\Gamma\left(y, y^{\prime}\right)$ is closed, tending $n$ to infinity, we have $P_{M}(, y) \otimes P_{M}\left(, y^{\prime}\right)\left(\Gamma\left(y, y^{\prime}\right)\right)=1$ for all $y, y^{\prime} \notin N$.

## §3. Proofs of the theorems

In this section we assume that the covariance operators of measures in $\mathscr{P}$ are non-singular since the other case is reduced easily to this case. For each $\mu$ $\in \mathscr{P}$ we denote by $\mathscr{F}(\mu)$ the set of all Borel measurable mappings $f$ from $E$ into itself with $\gamma_{\mu}\left(f^{-1}(A)\right)=\mu(A)$ for all $A \in \mathscr{E}$. For each $f \in \mathscr{F}(\mu)$ there exists exactly one probability measure $M_{f}$ in $\mathscr{M}(\mu)$ with $M_{f}(A \times B)=\gamma_{\mu}\left(f^{-1}(A) \cap B\right)$. We denote by $\mathscr{M}_{F}(\mu)$ the set of all $M_{f}$ with $f \in \mathscr{F}(\mu)$.

First we shall prove Theorem 1. To this end it is enough to show that
$\mathscr{M}_{\Gamma}(\mu) \subset \mathscr{M}_{F}(\mu)$, for, by Lemma 1 and Lemma $2, \mathscr{M}_{0}(\mu) \neq \varnothing$ and $\mathscr{M}_{0}(\mu) \subset \mathscr{M}_{F}(\mu)$. Now let $M \in \mathscr{M}_{\Gamma}(\mu)$. By Lemma 3, there exist a set $N \in \mathscr{E}$ and a regular conditional probability $P_{M}$ of $M$ with $\gamma_{\mu}(N)=0$ and $P_{M}(, y) \otimes P_{M}\left(, y^{\prime}\right)\left(\Gamma\left(y, y^{\prime}\right)\right)=1$ for all $y, y^{\prime} \notin N$. We denote by $S(y)$ the support of $P_{M}(, y)$; the smallest closed set with full measure. For each $y, y^{\prime} \notin N$, since $\Gamma\left(y, y^{\prime}\right)^{c}$ is open and $P_{M}(, y)$ $\otimes P_{M}\left(, y^{\prime}\right)\left(\Gamma\left(y, y^{\prime}\right) c\right)=0, S(y) \times S\left(y^{\prime}\right) \subset \Gamma\left(y, y^{\prime}\right)$, that is, $<x-x^{\prime}, y-y^{\prime}>\geqq 0$ for all $\left(x, x^{\prime}\right) \in S(y) \times S\left(y^{\prime}\right)$. Let $V$ be the covariance operator of $\mu$ and $\left(e_{n}\right)_{n \geqq 1}$ an orthonormal basis of $E$ consisting of eigen vectors of $V$. For each $n \geqq 1$, we set $\gamma_{n}=\gamma_{\mu} \operatorname{pr}_{\left[e_{n}\right]}^{-1}$ and $\gamma_{n}^{1}=\gamma_{\mu}\left(\operatorname{pr}_{\left[e_{n}\right]^{1}}\right)^{-1} .^{2}$ ) Notice that $\gamma_{n}$ is equivalent (mutually absolutely continuous) to the one-dimensional Lebesgue measure. A family $(C(\eta))_{\eta \in R^{1}}$ of non-empty subsets of $R^{1}$ is said to be monotone if there is a set $N_{0}$ with Lebesgue measure zero such that if $\eta, \eta^{\prime} \notin N_{0}$ then $\left(\xi-\xi^{\prime}\right)\left(\eta-\eta^{\prime}\right) \geqq 0$ for all $\left(\xi, \xi^{\prime}\right) \in C(\eta) \times C\left(\eta^{\prime}\right)$. It is known that if $(C(\eta))_{\eta \in \mathbb{R}^{1}}$ is monotone, $\operatorname{Card}(C(\eta))=1$ for all $\eta \notin N_{0}$ (see [2], the following statements are also the same as those of [2]). For each $z \in\left[e_{n}\right]^{\perp}$, we set

$$
\begin{aligned}
& S_{n}^{z}(\eta)=\left\{<x, e_{n}>: x \in S\left(\eta e_{n}+z\right)\right\}, \\
& N_{n}^{z}=\left\{\eta \in R^{1}: \eta e_{n}+z \in N\right\} .
\end{aligned}
$$

Assume that $\gamma_{n}\left(N_{n}^{z}\right)=0$ for some $z \in\left[e_{n}\right]^{\perp}$. Then the Lebesgue measure of $N_{n}^{z}$ is equal to zero and if $\eta, \eta^{\prime} \notin N_{n}^{z}$ then, for any $\left(\left\langle x, e_{n}\right\rangle,\left\langle x^{\prime}, e_{n}\right\rangle\right) \in S_{n}^{z}(\eta)$ $\times S_{n}^{z}\left(\eta^{\prime}\right)$,

$$
\begin{aligned}
& \left(<x, e_{n}>-<x^{\prime}, e_{n}>\right)\left(\eta-\eta^{\prime}\right) \\
& \quad=<x-x^{\prime},\left(\eta e_{n}+z\right)-\left(\eta^{\prime} e_{n}+z\right)>\geqq 0,
\end{aligned}
$$

which implies that $\left(S_{n}^{z}(\eta)\right)_{\eta \in R^{1}}$ is monotone. Thus Card $\left(S_{n}^{z}(\eta)\right)=1$ for all $\eta \notin N_{n}^{z}$. Now let

$$
\begin{aligned}
& B_{n}=\left\{z \in\left[e_{n}\right]^{\perp}: \gamma_{n}\left(N_{n}^{z}\right)=0\right\}, \\
& D_{n}=\left(\operatorname{pr}_{\left[e_{n}\right]^{1}}\right)^{-1}\left(B_{n}^{c}\right)
\end{aligned}
$$

and $D=U_{n \geqq 1} D_{n}$. Since $\gamma_{\mu}(N)=0$, using the Fubini theorem, we can easily prove that $\gamma_{\mu}\left(D_{n}\right)=\gamma_{n}^{1}\left(B_{n}^{c}\right)=0$ for all $n \geqq 1$ and hence, $\gamma_{\mu}(D)=0$. Now let $y \notin D \cup N$ and let $y=\eta_{n} e_{n}+z_{n}$ for each $n \geqq 1$, where $z_{n} \in\left[e_{n}\right]^{\perp}$. Then $\gamma_{n}\left(N_{n}^{z_{n}}\right)=0$ and $\eta_{n}$ $\notin N_{n}^{z_{n}}$ for all $n \geqq 1$, and so, $\left.\operatorname{Card}\left(\left\{<x, e_{n}\right\rangle: x \in S(y)\right\}\right)=1$ for all $n \geqq 1$, which implies that $\operatorname{Card}(S(y))=1$ for all $y \notin D \cup N$. We define a mapping $f$ of $E$ into itself by
2) For each $A \subset E,[A]$ denotes the closed linear subspace generated by $A$, and $\operatorname{pr}_{[A]}$ is the orthogonal projection to $[A]$. If $(X, \mathscr{X}, \nu)$ is a measure space and $f$ is measurable mapping of $X$ into $(Y, \mathscr{Y})$, we use $\nu f^{-1}$ to denote the measure defined by $\nu\left(f^{-1}(B)\right)(B \in \mathscr{Y})$.

$$
f(y)= \begin{cases}x_{y} & \text { for } y \notin D \cup N, \\ 0 & \text { for } y \in D \cup N,\end{cases}
$$

then for each $y \notin D \cup N, P_{M}(, y)=\varepsilon_{f(y)}$, the unit distribution at $f(y)$. Thus $M=M_{f}$ with $f \in \mathscr{F}(\mu)$.

Next we shall prove Theorem 2. Let $\mu_{1}, \mu_{2} \in \mathscr{P}$. For any $M_{i} \in \mathscr{M}\left(\mu_{i}\right)$ $(i=1,2)$, since $M_{1} * M_{2} \in \mathscr{M}\left(\mu_{1} * \mu_{2}\right)$,

$$
\begin{aligned}
e\left[\mu_{1} * \mu_{2}\right] & \leqq e\left[\mu_{1} * \mu_{2} ; M_{1} * M_{2}\right] \\
& =e\left[\mu_{1} ; M_{1}\right]+e\left[\mu_{2} ; M_{2}\right],
\end{aligned}
$$

it follows that $e\left[\mu_{1} * \mu_{2}\right] \leqq e\left[\mu_{1}\right]+e\left[\mu_{2}\right]$. Now assume that the equality holds. By Theorem 1, there exist $f_{i} \in \mathscr{F}\left(\mu_{i}\right)$ with $e\left[\mu_{i}\right]=e\left[\mu_{i} ; M_{f_{i}}\right](i=1,2)$. The relation;

$$
\begin{aligned}
e\left[\mu_{1} * \mu_{2} ; M_{f_{1}} * M_{f_{2}}\right] & =e\left[\mu_{1} ; M_{f_{1}}\right]+e\left[\mu_{2} ; M_{f_{2}}\right] \\
& =e\left[\mu_{1}\right]+e\left[\mu_{2}\right]=e\left[\mu_{1} * \mu_{2}\right]
\end{aligned}
$$

implies that $M_{f_{1}} * M_{f_{2}} \in \mathscr{M}_{0}\left(\mu_{1} * \mu_{2}\right)$ and hence, by Theorem 1 , there is an $f \in \mathscr{F}\left(\mu_{1} *\right.$ $\mu_{2}$ ) with $M_{f_{1}} * M_{f_{2}}=M_{f}$. Thus

$$
\begin{aligned}
& \gamma_{\mu_{1}} \otimes \gamma_{\mu_{2}}\left(\left\{\left(y, y^{\prime}\right): f_{1}(y)+f_{2}\left(y^{\prime}\right) \in A, y+y^{\prime} \in B\right\}\right) \\
& \quad=\gamma_{\mu_{1}} \otimes \gamma_{\mu_{2}}\left(\left\{\left(y, y^{\prime}\right): f\left(y+y^{\prime}\right) \in A, y+y^{\prime} \in B\right\}\right)
\end{aligned}
$$

for all $A, B \in \mathscr{E}$. Consequently $f\left(y+y^{\prime}\right)=f_{1}(y)+f_{2}\left(y^{\prime}\right)$ for almost all $\left(y, y^{\prime}\right)$ with respect to $\gamma_{\mu_{1}} \otimes \gamma_{\mu_{2}}$. Let $\left(e_{n}\right)_{n \geqq 1}$ be an orthonormal basis of $E$ consisting of eigen vectors of the covariance operator of $\mu_{1} * \mu_{2}$ and let $L_{n}=\left[e_{1}, e_{2}, \ldots, e_{n}\right]$. We denote by $\mathscr{L}_{n}$ the Borel $\sigma$-algebra of $L_{n}$ and set $\mathscr{B}_{n}=\operatorname{pr}_{L_{n}^{-1}}\left(\mathscr{L}_{n}\right)$. It is clear that $\left(\mathscr{B}_{n}\right)_{n \geqq 1}$ is increasing and generates $\mathscr{E}$. We denote by $\gamma^{(n)}$ and $\gamma_{i}^{(n)}$ [resp. $\sigma^{(n)}$ and $\left.\sigma_{i}^{(n)}\right]$ the probability measures on $L_{n}$ [resp. $\left.L_{n}^{\perp}\right]$ induced by $\mathrm{pr}_{L_{n}}$ [resp. $\left.\mathrm{pr}_{L_{n}^{\perp}}\right]$ from $\gamma$ and $\gamma_{\mu_{i}}$ respectively $(i=1,2)$, where $\gamma$ is a short for $\gamma_{\mu_{1}} * \gamma_{\mu_{2}}$. We consider the Bochner integrals;

$$
\begin{aligned}
f_{n}\left(y_{n}\right) & =\int_{L_{n}^{1}} f\left(y_{n}+z_{n}\right) d \sigma^{(n)}\left(z_{n}\right) \\
f_{n i}\left(y_{n}\right) & =\int_{L_{n}^{1}} f_{i}\left(y_{n}+z_{n}\right) d \sigma_{i}^{(n)}\left(z_{n}\right) \quad(i=1,2),
\end{aligned}
$$

where $y_{n} \in L_{n}$. Then $f_{n}\left(y_{n}+y_{n}^{\prime}\right)=f_{n 1}\left(y_{n}\right)+f_{n 2}\left(y_{n}^{\prime}\right)$ for almost all $\left(y_{n}, y_{n}^{\prime}\right)$ with respect to $\gamma_{1}^{(n)} \otimes \gamma_{2}^{(n)}$. Therefore $f_{n}$ is equal almost everywhere to an affine transformation from $L_{n}$ into $E$ and hence $\gamma f_{n}^{-1}$ is a Gaussian measure on $E$. On the other hand, since $\int_{E}\|f(y)\|^{2} d \gamma(y)<\infty$ and $E\left[f \mid \mathscr{B}_{n}\right]=f_{n}$ (as $E$-valued ran-
dom variables on $(E, \mathscr{E}, \gamma)$ ), by a martingale theorem, due to Chatterji [1], $\int_{E}\left\|f_{n}(y)-f(y)\right\|^{2} d \gamma(y) \rightarrow 0$ as $n \rightarrow \infty$. Therefore $\mu_{1} * \mu_{2}=\gamma f^{-1}$ is Gaussian and hence, $e\left[\mu_{1} * \mu_{2}\right]=0$. Consequently $e\left[\mu_{i}\right]=0(i=1,2)$, which implies both $\mu_{1}$ and $\mu_{2}$ are Gaussian. Thus the proof of Theorem 2 is completed.

## §4. An application

Let $(\Omega, \mathscr{B}, P)$ be a probability space. A measurable stochastic process $X(\omega)=(X(t, \omega))_{0 \leqq t<1}$ is considered as an $L_{2}[0,1)$-valued random variable if $E\left[\int_{0}^{1}(X(t))^{2} d t\right]<\infty$. We denote by $\mu_{X}$ the probability distribution (in $L_{2}[0,1)$ ) of $X$. In this section we use $e[X]$ to denote $e\left[\mu_{x}\right]$.

Now we consider a family of real random variables $\left\{\xi_{n j}: j=1,2, \ldots, 2^{n}-1\right.$, $n=1,2, \ldots\}$ with the following properties; (i) all $\xi_{n j}$ have the same distribution, (ii) for each $n \geqq 1,\left(\xi_{n j}\right)_{1 \leqq j<2^{n}}$ is independent and (iii) $E \xi_{n j}^{4}=c<\infty, E \xi_{n j}^{2}=1$ and $E \xi_{n j}=0$. Using the family we define a sequence of processes $\left(X_{n}\right)_{n \geqq 1}$ by

$$
X_{n}(t, \omega)=2^{-n / 2} \sum_{1 \leqq j<2^{n}} S_{n j}(\omega) f_{n j}(t)
$$

for $(t, \omega) \in[0,1) \times \Omega$, where $S_{n j}=\Sigma_{1 \leqq i \leqq j} \xi_{n i}$ and $f_{n j}$ are the indicators of intervals $\left[j 2^{-n},(j+1) 2^{-n}\right)$. It is known that the sequence $\left(X_{n}\right)_{n \geqq 1}$ of stochastic processes converges in law to a Brownian motion. The purpose of this section is to prove this in a wider space $L_{2}[0,1)$ by making use of the functional $e$. Since $E\left[\int_{0}^{1}\left(X_{n}(t)\right)^{2} d t\right]$ $=\left(1-2^{-n}\right) / 2<\infty, X_{n}$ are $L_{2}[0,1)$-valued random variables, and $\mu_{n}=\mu_{X_{n}}$ have the mean vectors 0 and the covariance operators $V_{n}$, the integral operators with kernel $v_{n}(s, t)=\left(\left[s 2^{-n}\right] \wedge\left[t 2^{-n}\right]\right) / 2^{n}$. We prove that $\left(\mu_{n}\right)_{n \geqq 1}$ converges to the Gaussian measure with mean vector 0 and covariance operator $V$, with kernel $v(s, t)=s \wedge t$. Using the random variables $S_{n j}^{0}=\sum_{1 \leqq i \leqq j} \xi_{n+1,2 i}$ and $S_{n j}^{1}=\sum_{1 \leqq i \leqq j} \xi_{n+1,2 i-1}$ ( $1 \leqq j<2^{n}$ ), we define the processes $X_{n}^{0}, X_{n}^{1}$ and $Z_{n}$ by

$$
\begin{aligned}
& X_{n}^{0}=2^{-n / 2} \sum_{1 \leqq j<2 n} S_{n j}^{0} f_{n j}, \\
& X_{n}^{1}=2^{-n / 2} \sum_{1 \leqq j<2^{n}} S_{n j}^{1} \tilde{f}_{n j} \quad\left(\tilde{f}_{n j}=f_{n+1,2 j-1}+f_{n+1,2 j}\right), \\
& Z_{n}=2^{-(n+1) / 2} \xi_{n+1,2^{n+1}-1} f_{n+1,2^{2+1-1}},
\end{aligned}
$$

respectively. By our assumptions, for each $n,\left(X_{n}^{0}, X_{n}^{1}, Z_{n}\right)$ is independent and $X_{n}, X_{n}^{0}$ and $X_{n}^{1}$ have the same distribution as random variables taking values in the $\left(2^{n}-1\right)$-dimensional Euclidean space. Since $X_{n+1}=2^{-1 / 2}\left(X_{n}^{0}+X_{n}^{1}\right)+Z_{n}$ and $e\left[Z_{n}\right] \leqq 4^{-n}$, we have $e\left[X_{n+1}\right] \leqq e\left[X_{n}\right]+4^{-n}$ for all $n \geqq 1$. Therefore the limit $\alpha=\lim _{n} e\left[X_{n}\right]$ exists. Let $g_{0}=1$ and $g_{m k}=2^{(m-1) / 2}\left(f_{m, k-1}-f_{m k}\right)$ for $k=1,3, \ldots$, $2^{m}-1$ and $m=1,2, \ldots$. It is known that $\left\{g_{0}, g_{m k}\right.$ : odd $\left.k<2^{m}, m \geqq 1\right\}$ is an orthonormal basis in $L_{2}[0,1)$ (the Haar functions). For any $m_{0}$, since

$$
\sup _{n} \sum_{\substack{\text { odd } k<2^{m} \\ m \geqq m_{0}}} E\left[<X_{n}, g_{m k}>^{2}\right] \leqq 2^{-\left(m_{0}-1\right)}
$$

$\left(\mu_{n}\right)_{n \geqq 1}$ is relatively, weakly compact [3; p. 154]. Hence we can choose a subsequence $\left(\mu_{n_{i}}\right)_{i \geqq 1}$ of $\left(\mu_{n}\right)_{n \geqq 1}$ that converges weakly to a probability measure $\mu$ on $L_{2}\left[0,1\right.$ ). From (iii) it follows that $\int\|x\|^{4} d \mu_{n}(x)=E\left[\left\|X_{n}\right\|^{4}\right] \leqq 2^{-n} c+3\left(1-2^{-n}\right)$ and $\int\|x\|^{4} d \gamma_{\mu_{n}}(x) \leqq 3$, and hence we have

$$
\begin{aligned}
& \lim _{m} \sup _{n} \int_{\{\|X\| \geqq m\}}\|x\|^{2} d \mu_{n}(x) \\
= & \lim _{m} \sup _{n} \int_{\{\|X\| \geqq m\}}\|x\|^{2} d \gamma_{\mu_{n}}(x)=0 .
\end{aligned}
$$

Using the relations, we can prove that $e[\mu]=\lim _{i} e\left[X_{n_{i}}\right] . \quad$ Let $\mu_{0}$ and $\mu_{1}$ be the limit distributions of $\left(X_{n_{i}}^{0} / \sqrt{2}\right)_{i \geqq 1}$ and $\left(X_{n_{i}}^{1} / \sqrt{2}\right)_{i \geqq 1}$ respectively. Then $\mu=\mu_{0} *$ $\mu_{1}$. On the other hand, since

$$
\begin{aligned}
\alpha & =e[\mu] \leqq\left(e\left[\mu_{0}\right]+e\left[\mu_{1}\right]\right) / 2 \\
& =\lim _{i}\left(e\left[X_{n_{i}}^{0} / \sqrt{2}\right]+e\left[X_{n_{i}}^{1} / \sqrt{2}+Z_{n}\right]\right) \\
& \leqq \lim _{i}\left(e\left[X_{n_{i}}^{0} / \sqrt{2}\right]+e\left[X_{n_{i}}^{1} / \sqrt{2}\right]+e\left[Z_{n_{i}}\right]\right) \\
& =\alpha / 2+\alpha / 2=\alpha
\end{aligned}
$$

we have $e[\mu]=e\left[\mu_{0}\right]+e\left[\mu_{1}\right]$, which implies, by Theorem 2 , that $\mu$ is Gaussian.

## References

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[^0]:    1) Throughout this paper, $\pi_{1}$ and $\pi_{2}$ denote the first and the second coordinate mappings of $E^{2}$ onto $E$.
