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# Albert Badrikian <br> Measurable linear mappings from a Wiener space 

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## $\mathcal{N u m b a m}^{\prime}$

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## Measurable linear mappings from a Wiener space

## CHAPTER ONE

## Gaussian probabilities

We want to describe a situation involving the Wiener measure and the white noise measure and we shall work with a structure more general than the "abstract Wiener space". More precisely instead of handling only separable Banach space, we shall work with Lusin spaces. It is not the maximum of generality, but we give only the preliminaries needed for our purpose.

## 1-General definitions and results

Let for a moment $E$ be a topological vector space (on $\mathbb{R}$ if it is not precised), locally convex and Hausdorff and let us denote $E^{\prime}$ its dual.

As it is well known a Borelian probability $P$ on $E$ is said "Gaussian centered" if for every $x^{\prime} \in E^{\prime}, x^{\prime}(P)$ is a Gaussian centered probability on $\mathbb{R}$ (eventually equal to $\delta_{0}$ ). This is equivalent to say that for every $x^{\prime} \in E^{\prime}$, the random variable : $\left\langle x^{\prime}, \bullet\right\rangle_{E^{\prime}, E}$ on $(E, P)$ is Gaussian centered.
In the following we shall omit the mention "centered".
If $P$ is a Gaussian probability on $E$, we associate to it a linear map $j$ from $E^{\prime}$ into $L^{2}(E, P)$ defined as :

$$
j\left(x^{\prime}\right)=<x^{\prime}, \bullet>_{E^{\prime}, E}
$$

(we have used the same notation for a map and its $P$-equivalence class, as we shall do often in the sequel).
We have not equipped $E^{\prime}$ with a topology, so we cannot speak about the continuity of $j$. The (algebraic) transpose of $j$ defines a linear map from $L^{2}(E, P)$ into the algebraic dual $\left(E^{\prime}\right)^{*}$ of $E^{\prime}$. We shall denote by $S$ this transpose.

We have therefore :

$$
\text { for } \begin{aligned}
f \in L^{2}(E, P): \quad\left(j\left(x^{\prime}\right), f\right)_{L^{2}(E, P)} & =\int_{E}<x^{\prime}, x>f(x) P(d x) \\
& =<S f, x^{\prime}>_{\left(E^{\prime}\right)^{*}, E^{\prime}}
\end{aligned}
$$

Now we can give another definition for $S$ :
for $f \in L^{2}(E, P)$, let us consider the vector function : $x \rightsquigarrow f(x) x$. It belongs scalarly in $L^{1}(E, P)$, that is to say,
for every $\mathrm{x}^{\prime} \in E^{\prime}: x \rightsquigarrow f(x)<x^{\prime}, x>\quad$ is integrable with respect to P .

Then we define its weak integral as the unique element from $\left(E^{\prime}\right)^{*}$ such that,

$$
\text { for every } \mathrm{x}^{\prime} \in E^{\prime}:<\int f(x) x P(d x), x^{\prime}>_{E^{\prime *}, E^{\prime}}:=\int_{E}<x^{\prime}, x>f(x) P(d x) .
$$

Now $S f$ is the barycenter of the measure $f(x) P(d x)$.
We shall denote by $H(P)$ (or simply $H$ if there is no ambiguity) the subspace $S\left(L^{2}(E, P)\right.$ ) of $\left(E^{\prime}\right)^{*}$.
$H(P)$ is called the "Cameron-Martin space" of $P$.
The closure of $j\left(E^{\prime}\right)$ in $L^{2}(E, P)$ is called the "Gaussian space" of $P$ (or the "first
Wiener chaos") and will be denoted by $E_{2}^{\prime}(P)$. It is an Hilbert space for the topology induced by the topology of $L^{2}(E, P)$.

Let us notice that :

- $H(P)=S\left(E_{2}^{\prime}(P)\right)$ and that :
- $S$, when restricted to $E_{2}^{\prime}(P)$, is injective, whence bijective, with values in $H(P)$.

Actually :
Suppose first that $f \in L^{2}(E, P)$ is orthogonal to $E_{2}^{\prime}(P)$; that means that :

$$
\text { for every } x^{\prime} \in E^{\prime}: \int_{E} j\left(x^{\prime}\right)(x) f(x) P(d x)=0
$$

Then we have

$$
\begin{aligned}
\text { for every } x^{\prime} \in E^{\prime} & : \int_{E}<x^{\prime}, x>f(x) P(d x) \\
& =<\int_{E} f(x) x P(d x), x^{\prime}>=0
\end{aligned}
$$

But this means that :

$$
S f=0
$$

The first assertion is therefore proven.

For proving the second assertion it suffices to notice that if

$$
f \in E^{\prime 2}(P)
$$

is such that

$$
S(f)=0,
$$

then :

$$
f \perp E_{2}^{\prime}(P) \quad \text { and } \quad f \perp f, \quad \text { whence } f=0
$$

Now we can put on $H(P)$ an Hilbertian structure : the image of the Hilbertian structure of $E_{2}^{\prime}(P)$ by the linear isomorphism $S$.
Then we have :
a) $(S f, S g)_{H(P)}=\int_{E} f(x) g(x) P(d x), \quad \forall(f, g) \in\left(E_{2}^{\prime}(P)\right)^{2}$,
b) if $f \in E_{2}^{\prime}(P)$ and $y^{\prime} \in E^{\prime}:$

$$
<S f, y^{\prime}>_{\left(E^{\prime}\right)^{*}, E^{\prime}}=\int_{E} f(x)<y^{\prime}, x>P(d x)=\left(f, j\left(y^{\prime}\right)\right)_{E_{2}^{\prime}(P)}
$$

c) if $\left(x^{\prime}, y^{\prime}\right) \in E^{\prime} \times E^{\prime}$ :

$$
\left(S \circ j\left(x^{\prime}\right), S \circ j\left(y^{\prime}\right)\right)_{H(P)}=\int_{E}<x^{\prime}, x><y^{\prime}, x>P(d x) .
$$

Remark : If $\operatorname{Supp} P=E$, the map $j: E^{\prime} \rightarrow E_{2}^{\prime}(P)$ is injective. Actually, if $j\left(x^{\prime}\right)=0$, then :

$$
\int<x^{\prime}, x>^{2} P(d x)=0
$$

therefore :

$$
\left.<x^{\prime}, \bullet\right\rangle \quad \text { is null, almost everywhere. }
$$

But $\left\langle x^{\prime}\right.$, . $>$ being continuous, we have :

$$
\left\langle x^{\prime}, \bullet\right\rangle=0 \quad \text { everywhere } .
$$

For the questions about which we shall be concerned, the case where $H(P) \subset E$ is of paramount importance . We shall give a sufficient condition under which $H(P) \subset E$.

If $h \in H(P)$ we shall denote by $\tilde{h}$ the element $S^{-1}(h)$ of $E_{2}^{\prime}(P)$.

Lemma 1: Let us suppose $H(P) \subset E$. If $h \in H(P)$, let us denote by $P_{h}$ the translate of $P$ by $h,\left(P_{h}(A)=P(A+h)\right)$. Then we have :

- $P_{h} \equiv P$
- $\frac{d P_{h}}{d P}=\exp \left\{\tilde{h}-\frac{1}{2}\|h\|_{H(P)}^{2}\right\}$.


## Proof :

( We shortly write $H$ instead of $H(P)$ ). $\widetilde{h}$ being a ( $P$-class of ) Gaussian random variable with variance $\|h\|_{H}^{2}$, we have:

$$
\int_{E} \exp (\widetilde{h}) d P=\exp \left\{\frac{1}{2}\|h\|_{H}^{2}\right\}
$$

Therefore :

$$
\exp \left\{\tilde{h}-\frac{1}{2}\|h\|_{H}^{2}\right\} \cdot P \quad \text { is a probability. }
$$

Now, let us recall some results about the Fourier transform (or characteristic function) of a Borelian probability on E , not necessarily Gaussian centered :

If $Q$ is a Borelian probability on $E$, we define its characteristic function as the map : $E^{\prime} \rightarrow \mathbb{C}$ defined by :

$$
x^{\prime} \rightsquigarrow \int e^{i\left\langle x^{\prime}, x\right\rangle_{E^{\prime}, E}} Q(d x):=\widehat{Q}\left(x^{\prime}\right) .
$$

Then two probabilities coincide if and only if their characteristic functions coincide.
In our case, where $P$ is Gaussian, we have :

$$
\widehat{P}\left(x^{\prime}\right)=\exp \left\{-\frac{1}{2}\left\|j\left(x^{\prime}\right)\right\|_{L^{2}(E, P)}^{2}\right\}=\exp \left\{-\frac{1}{2}\left\|S j\left(x^{\prime}\right)\right\|_{H(P)}^{2}\right\}
$$

Now :

$$
\widehat{P}_{h}\left(x^{\prime}\right)=\exp \left\{i<x^{\prime}, h>\right\} \widehat{P}\left(x^{\prime}\right)
$$

Moreover we have :

$$
\begin{aligned}
& \int_{E} \exp \left\{i<x^{\prime}, x>\right\} \exp \left\{\tilde{h}(x)-\frac{1}{2}\|h\|_{H}^{2}\right\} P(d x) \\
& =\exp \left\{-\frac{1}{2}\|h\|_{H}^{2}\right\} \int \exp \left\{i\left(<x^{\prime}, x>-i \widetilde{h}(x)\right)\right\} P(d x) \\
& =\exp \left\{-\frac{1}{2}\|h\|_{H}^{2}\right\} \exp \left\{\frac{1}{2}\|h\|_{H}^{2}-\frac{1}{2}\left\|j x^{\prime}\right\|_{E_{2}^{\prime}(P)}^{2}+2 i\left(j\left(x^{\prime}\right), \tilde{h}\right)_{L^{2}(E, P)}\right\}
\end{aligned}
$$

Actually, the 2-dimension random vector $\left(\widetilde{h}, j\left(x^{\prime}\right)\right)$ is Gaussian (centered) and its covariance matrix is equal to :

$$
\left(\begin{array}{cc}
\left\|j\left(x^{\prime}\right)\right\|_{E_{2}^{\prime}}^{2} & \left(j\left(x^{\prime}\right), \widetilde{h}\right)_{E_{2}^{\prime}} \\
\left(j\left(x^{\prime}\right), \widetilde{h}\right)_{E_{2}^{\prime}} & \|h\|_{H}^{2}
\end{array}\right)
$$

From what precedes we deduce immediately that $P_{h}$ and $\exp \left\{\tilde{h}-\frac{1}{2}\|h\|_{H}^{2}\right\} P$ have the same characteristic functional.

- Lemma 1 is proven.-

Corollary : If $F$ is a $P$-measurable vector subspace, carrying $P$, then : $H(P) \subset F$.
(We suppose that $H(P) \subset E$ ).
Proof :
Let us suppose the contrary : there exists $h \in H(P)$ such that $h \notin F$. Then :

$$
(h+F) \cap F=\emptyset
$$

But

$$
P(h+F)>0, \quad \text { by Lemma } 1
$$

and since : $P(F)=1$, we have :

$$
P((h+F) \cap F)>0 .
$$

There is a contradiction.

- The Corollary is proven. -

Remark : we shall see later that for every $P$-measurable vector space we have : $P(F)=0$ or 1 (we shall give the proof if $E$ is Lusin).

Lemma 2 : Let $E$ and $F$ be two Hausdorff l.c.v.s., and $P_{1}$ be a Gaussian probability on $E$. Let $u: E \rightarrow F$ continuous linear and $P_{2}=u\left(P_{1}\right)$, then $P_{2}$ is Gaussian. Let $H_{i}:=H\left(P_{i}\right),(i=1,2)$ then :

- if $H_{1} \subset E$ we have : $H_{2} \subset F$.


## Moreover :

- $u\left(H_{1}\right)=H_{2}$,
- $u_{\mid H_{i}}$ is continuous from $H_{1}$ into $H_{2}$ (for the Hilbertian topologies).
- the image of the unit ball of $H_{1}$ by $u^{\prime *}$ is the unit ball of $H_{2}$.


## Proof :

The map $u$ induces a linear map $U$ from $L^{2}\left(F, P_{2}\right)$ into $L^{2}\left(E, P_{1}\right)$ defined by :

$$
U f=f \circ u
$$

It is easy to see that $U$ is an isometry from $F_{2}^{\prime}\left(P_{2}\right)$ into $E_{2}^{\prime}\left(P_{1}\right)$; let $G=U\left(F_{2}^{\prime}\left(P_{2}\right)\right)$.
Let $S_{1}$ (resp. $S_{2}$ ) the barycenter map considered as a map from $E_{2}^{\prime}\left(P_{1}\right)\left(\operatorname{resp} . F_{2}^{\prime}\left(P_{2}\right)\right)$ onto $H_{1}$ (resp. $H_{2}$ ). $S_{i}$ is an isomorphism between Hilbert spaces.
I assert that :

$$
u^{\prime *} \circ S_{1} \circ U=S_{2}
$$

In fact, let $g \in F_{2}^{\prime}\left(P_{2}\right)$ then :

$$
\begin{aligned}
u^{\prime *} \circ S_{1} \circ U(g) & =u^{\prime *}\left(\int_{E} g(u(x)) x P_{1}(d x)\right) \\
& =\int_{E} g(u(x)) u(x) P_{1}(d x) \\
& =\int_{F} g(y) y P_{2}(d y) \\
& =S_{2}(g)
\end{aligned}
$$

We have proven that :

$$
u^{\prime *}\left(H_{1}\right) \supset u^{\prime *}\left(S_{1}(G)\right)=H_{2}
$$

Moreover, if $h \in H_{1}$ with $\tilde{h}$ orthogonal to $G$, we have

$$
u^{\prime *}(h)=0, \text { so } u^{\prime *}\left(H_{1}\right)=H_{2}
$$

In fact,

$$
\text { for each } y^{\prime} \in F^{\prime}:<u^{\prime *}\left(S_{1} h\right), y^{\prime}>=u^{\prime}\left(y^{\prime}\right) \int_{E} \widetilde{h}(x) x P_{1}(d x), ~ \begin{aligned}
& =\int_{E} \widetilde{h}(x) y^{\prime}(u(x)) P_{1},(d x) \\
& =0
\end{aligned}
$$

Now it is clear that
$u_{\mid H_{1}}^{\prime *}$ is continuous from $H_{1}$ to $H_{2}$, with norm one.

- Lemma 2 is proven. -

In the sequel we shall do the following hypothesis about $P$ :
(H) for every $\varepsilon>0$, there exists a compact $K_{\varepsilon} \subset E$ such that :

$$
P\left(K_{\varepsilon}\right)>1-\varepsilon
$$

A probability (not necessarily Gaussian) satisfying the hypothesis $(H)$ is said a "Radon probability". A Borelian probability will be, ipso facto, a Radon probability in the following cases (that we shall consider in the sequel) :

- $E$ is a separable Banach space
- more generally : $E$ is a Lusin space.

Lemma 3 : If $P$ is a Radon probability (not necessarily Gaussian), $\widehat{P}$ is continuous from $E_{c}^{\prime}\left(E^{\prime}\right.$ with the topology of compact convergence) into $\mathbb{C}$.

## Proof :

Let $\varepsilon \in] 0,1\left[\right.$ be given and let $K_{\varepsilon}$ be a compact subset of $E$ such that :

$$
P\left(K_{\varepsilon}\right)>1-\varepsilon
$$

Let at last $p_{K_{\varepsilon}^{\circ}}$ the gauge of the absolute polar of $K_{\varepsilon}$, that is to say :

$$
\begin{gathered}
K_{\varepsilon}^{\circ}=\left\{x^{\prime} \in E^{\prime}, \quad\left|<x^{\prime}, x>\right| \leq 1, \quad \forall x \in K_{\varepsilon}\right\} \\
p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right)=\inf \left\{\lambda, \quad \lambda>0, \quad x^{\prime} \in \lambda K_{\varepsilon}^{\circ}\right\}
\end{gathered}
$$

I assert that :

$$
\left|1-\widehat{P}\left(x^{\prime}\right)\right| \leq 2 \varepsilon+p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right), \quad \forall x^{\prime}
$$

Actually, on one hand we have :

$$
x^{\prime}\left(K_{\varepsilon}\right) \subset\left[-p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right), p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right)\right], \quad \text { by definition }
$$

and on the other hand, if we denote by $\nu$ the probability $x^{\prime}(P)$ :

$$
\begin{aligned}
\left|1-\widehat{P}\left(x^{\prime}\right)\right| & =\left|\int_{\mathbb{R}}\left(1-e^{i t}\right) \nu(d t)\right| \\
& \leq \int_{\mathbb{R}}\left|1-e^{i t}\right| \nu(d t)
\end{aligned}
$$

But :

$$
\int_{\mathbb{R}}=\int_{|t| \leq p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right)}+\int_{|t|>p_{K_{\varepsilon}^{\circ}}^{\circ}\left(x^{\prime}\right)}
$$

and

$$
\left|1-e^{i t}\right| \leq \inf (2,|t|), \quad \forall t \in R .
$$

Therefore :

$$
\int_{|t| \leq p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right)}\left|1-e^{i t}\right| d \nu(t) \leq p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right)
$$

and :

$$
\int_{|t|>p_{K_{\varepsilon}^{\circ}}\left(x^{\prime}\right)}\left|1-e^{i t}\right| d \nu(t) \leq 2 \varepsilon
$$

and the assertion is proven.
We deduce then immediately that :

$$
x^{\prime} \in \varepsilon K_{\varepsilon}^{\circ} \Rightarrow\left|1-\widehat{P}\left(x^{\prime}\right)\right| \leq 3 \varepsilon,
$$

- Lemma 3 is proven. -

As a corollary we deduce that, if $P$ is Gaussian, the map $j: E^{\prime} \rightarrow L^{2}(E, P)$ is continuous for $E_{c}^{\prime}$.

Lemma 4 : Let $P$ be a Gaussian (Radon) probability on $E$. If $E$ is quasi-complete (this means that the closed disked hull of every bounded subset is complete), then we have: $H(P) \subset E$.

## Proof :

We have noticed that :

$$
j: E_{c}^{\prime} \rightarrow L^{2}(E, P) \quad \text { is continuous. }
$$

Therefore, for every $f \in L^{2}(E, P)$, the linear form :

$$
x^{\prime} \rightsquigarrow \int_{E} f(x)<x^{\prime}, x>P(d x) \quad \text { is continuous on } E_{c}^{\prime} .
$$

By Mackey's theorem, we deduce that :

$$
\int f(x) x P(d x) \in E
$$

Remark : a complete space being trivially quasi-complete, Lemma 4 can be applied when $E$ is complete.

Now we have the main following result :
THEOREM 1 : If $P$ is Radon, Gaussian, then $H(P) \subset E$.

## Proof :

Let $\tilde{E}$ be the completion of $E$ and $i$ be the canonical injection $E \rightarrow \widetilde{E}$. Let $\widetilde{P}=i(P)$. By Lemma 2 and Lemma 4, we have :

$$
i(H(P))=H(\widetilde{P}) \subset \widetilde{E}
$$

But $E=i(E)$ is a subspace carrying $\tilde{P}$. Therefore, by corollary of Lemma 1 :

$$
H(\widetilde{P}) \subset E
$$

It is easy to see that $H(P)$ and $H(\widetilde{P})$ are isometric.
Remark : under the hypothesis of the theorem, the canonical injection $H(P) \rightarrow E$ is continuous.

For proving this fact it suffices to prove that, for every $f \in L^{2}(E, P)$ and every continuous semi-norm $q$, we have :

$$
q\left(\int f(x) x P(d x)\right) \leq c\|f\|_{L^{2}(E, P)}, \quad \text { where } c \text { is a constant. }
$$

By a result of Fernique and Skorokhod, we know that, for every continuous semi-norm $q$ on $E$, we have :

$$
\int_{E}|q(x)|^{2} P(d x)<\infty
$$

(A non-trivial fact if $E$ is infinite-dimensional).
Therefore :

$$
\begin{aligned}
q\left(\int_{E}(x f(x) P(d x))\right. & \leq \int|f(x)| q(x) P(d x) \\
& \leq\|f\|_{L_{(E, P)}^{2}}\|q\|_{L_{(E, P)}^{2}}
\end{aligned}
$$

The assertion is proven.

We shall call a "Gauss space" (or "abstract Wiener space") a couple ( $E, P$ ) where $E$ is a Hausdorff l.c.s. and $P$ is a centered Gaussian Radon measure on $E$.

## 2 - Cameron-Martin space of a Lusin Gauss space

From now on, $E$ will denote a Lusin locally convex vector space and $P$ a Gaussian Borelian probability on $E$ ( $P$ is automatically Radon).

We shall suppose that $\operatorname{Supp} P=E$. (This is not a loss of generality). We shall suppose moreover (unless the contrary is specified) that $\operatorname{dim} E=+\infty$. (The case where $\operatorname{dim} E<\infty$ being well known).

As a first example of this situation we have the "discrete white noise" or the canonical Gaussian measure on $\mathbb{R}^{\mathbb{N}}$ which is described as follows :

Let $E=\mathbb{R}^{\mathbb{N}}$ with the product topology (it is a polish space, therefore Lusinian). Its dual $E^{\prime}$ is the space $\mathbb{R}_{0}^{\mathbb{N}}$ of sequences of real numbers "almost null". Let $P=\gamma_{1}^{\otimes \mathbb{N}}$ the canonical Gaussian measure on $E$ ( $\gamma_{1}$ is the $\mathcal{N}(0,1)$ probability on $R$ ).

Here we have $E^{\prime} \subset E$. If $x^{\prime}=\left(x_{n}^{\prime}\right) \in E^{\prime}, j\left(x^{\prime}\right)$ is the ( $P$-class) function $x \rightsquigarrow \sum x_{n} x_{n}^{\prime}$ where $\left.x=\left(x_{n}\right)_{n \in \mathbb{N}} \in \mathbb{R}^{\mathbb{N}}\right)$.
Obviously :

$$
\left\|j\left(x^{\prime}\right)\right\|_{L^{2}\left(\mathbb{R}^{\mathbb{N}}, P\right)}^{2}=\sum\left|x_{n}^{\prime}\right|^{2}=\left\|x^{\prime}\right\|_{\ell^{2}}^{2} .
$$

Therefore

$$
E_{2}^{\prime}(P) \simeq \ell^{2}
$$

Let us now determine the Cameron-Martin space of $P$.
Let $x^{\prime} \in E^{\prime}$, then

$$
S\left(j\left(x^{\prime}\right)\right)=\int_{\mathbb{R}^{\mathbb{N}}}\left(\sum_{n} x_{n}^{\prime} x_{n}\right) x P(d x)
$$

Therefore :

$$
S\left(j\left(x^{\prime}\right)\right)_{k}=\int_{\mathbb{R}^{\mathbb{N}}}\left(\sum_{n} x_{n}^{\prime} x_{n}\right) x_{k} P(d x)=x_{k}^{\prime}: \quad \forall k
$$

We have then

$$
S\left(j\left(E^{\prime}\right)\right)=E^{\prime} \quad \text { and } \quad H(P)=\ell^{2}
$$

The above example is actually the model for a Gaussian probability on a Lusin space, as we shall now see.

THEOREM 2 : Let $E$ be a Lusin l.c.s. (Hausdorff) with infinite dimension and let $P$ be a Gaussian probability on $E$, whose support is equal to $E$. Then there exists a continuous linear injection from $E$ into $\mathbb{R}^{\mathbb{N}}$, denoted $\pi$, such that $\pi(P)$ is the canonical Gaussian measure on $\mathbb{R}^{\mathbb{N}}$.

## Proof :

The family of continuous functions on $E$ :

$$
\left\{\left\langle x^{\prime}, \cdot\right\rangle_{E^{\prime}, E}, x^{\prime} \in E^{\prime}\right\} \quad \text { is separating. }
$$

$E$ being Lusin and $\operatorname{dim} E=+\infty$ there exists an infinite sequence $\left(x_{n}^{\prime}\right)_{n \in \mathbb{N}}$ in $E^{\prime}$ separating points of $E$. Then $\sigma\left\{x_{n}^{\prime}, n \in \mathbb{N}\right\}$ is identical to the Borelian $\sigma$-field of $E$, and the $j\left(x_{n}^{\prime}\right)$ generate $E_{2}^{\prime}(P)$.
Let $\left(y_{n}^{\prime}\right)$ be the orthonormal sequence obtained from the $\left(x_{n}^{\prime}\right)$ by application of the Schmidt orthogonalization procedure to $\left(j\left(x_{n}^{\prime}\right)\right)_{n}$. The $\left(y_{n}^{\prime}\right)$ separate points of $E$, and

$$
\operatorname{card}\left\{y_{n}^{\prime}, n \in \mathbb{N}\right\}=+\infty
$$

Let $\pi$ be the map $x \rightsquigarrow\left(\left\langle y_{n}^{\prime}, x\right\rangle\right)_{n}$ from $E$ into $\mathbb{R}^{\mathbb{N}}$. It is trivial to see that $\pi$ is linear, continuous, injective and that :

$$
\pi(P)=\gamma_{1}^{\otimes \mathbf{N}}
$$

-Q.E.D. -

Remark 1 : $\pi$ being continuous, it is clearly Borelian. Moreover $\pi$ being injective, $E$ and $\mathbb{R}^{\mathbb{N}}$ being Lusinian, $\pi$ is bimeasurable (not necessarily onto).

Remark 2: It is easy to see that if $E$ is finite dimensional (or more generally $H(P)$ is finite dimensional) there exists $n$ and a linear continuous (not necessarily injective) application $\pi: E \rightarrow \mathbb{R}^{n}$, such that $\pi(P)=\gamma_{n}$.

Coming back to the situation $\operatorname{dim} E=+\infty, \operatorname{Supp} P=E$, we see that $H(P)$ is separable, therefore a Polish space. Therefore, $H(P)$ is a Borel subset of $E$, and we can speak of $P(H(P))$.

Proposition 1 : If $H$ denotes the Cameron-Martin space of $E$, then $P(H)=0$.
Proof :
We can suppose that $(E, P)$ is the canonical Gauss space $\mathbb{R}^{\mathbb{N}}$. Then :

$$
H=\ell^{2}
$$

Now the random variables $x \rightsquigarrow x_{n}$ on $\mathbb{R}^{\mathbb{N}}$ are independent and $\mathcal{N}(0,1)$. Therefore :

$$
\sum x_{n}^{2}=+\infty, \quad \text { almost surely }
$$

or equivalently :

$$
P\left(\ell^{2}\right)=0
$$

- Q.E.D. -

Remark : The result is clearly false if $\operatorname{dim} H<\infty$. For instance, $\gamma_{n}\left(\mathbb{R}^{n}\right)=1$ (in this case $H=\mathbb{R}^{n}$ ).

We can actually prove that $P(H)=0$ if and only if $\operatorname{dim} H=\infty$. Let us notice that in any case $P(H)=0$ or 1 .

Proposition 2: If $\left(\widetilde{x}_{n}\right)$ is an orthonormal basis of $E_{2}^{\prime}(P)$ and $x_{n}=S\left(\widetilde{x}_{n}\right)$, then for every $x^{\prime} \in E^{\prime}:\left(\sum_{n \leq k} \tilde{x}_{n}(\cdot)\left\langle x^{\prime}, x_{n}\right\rangle_{E^{\prime}, E}\right)_{k} \longrightarrow\left\langle x^{\prime}, \bullet\right\rangle$, almost surely.

## Proof :

The random variables $\tilde{x}_{n}(\cdot)\left\langle x^{\prime}, x_{n}\right\rangle_{E^{\prime}, E},(n \in \mathbb{N})$ are independent Gaussian with variances

$$
\left\langle x^{\prime}, x_{n}\right\rangle_{E^{\prime}, E}^{2}
$$

But

$$
\left\langle x^{\prime}, x_{n}\right\rangle_{E^{\prime}, E}=\left(S j\left(x^{\prime}\right), x_{n}\right)
$$

Therefore :

$$
\sum_{n}\left|\left\langle x^{\prime}, x_{n}\right\rangle\right|^{2}=\left\|S j\left(x^{\prime}\right)\right\|_{H}^{2}<\infty
$$

and

$$
\sum_{n} \tilde{x}_{n}(\cdot)\left\langle x^{\prime}, x_{n}\right\rangle \text { converges almost surely. }
$$

We have to prove that the limit is $j\left(x^{\prime}\right)$.
Actually we have

$$
\begin{aligned}
& \mathbb{E}\left\{j\left(x^{\prime}\right) \mid \widetilde{x}_{i}, \quad 0 \leq i \leq n\right\}=\sum_{i=0}^{n} \mathbb{E}\left\{j\left(x^{\prime}\right) \widetilde{x}_{i}\right\} \widetilde{x}_{i} \\
& =\sum_{i=0}^{n}\left(S j\left(x^{\prime}\right), x_{i}\right) \widetilde{x}_{i}=\sum_{i=0}^{n}\left\langle x^{\prime}, x_{i}\right\rangle_{E^{\prime}, E} \widetilde{x}_{i} .
\end{aligned}
$$

Now it suffices to apply the martingale convergence theorem (for real martingales).

- Q.E.D. -

Remark : The preceding result is very poor. We did not prove that

$$
\sum_{n \leq k} \widetilde{x}_{n} \cdot x_{n} \longrightarrow I d \text { for the weak topology, }
$$

the exceptional set of non convergence being a priori dependent of $x^{\prime}$.

In chapter two we shall prove that $\sum_{n \leq k} \tilde{x}_{n} \cdot x_{n} \longrightarrow I d$, for the topology of $E$.
Proposition 3 : The canonical injection $i: H(P) \rightarrow E$ is compact. We have even more if $K$ denotes the unit ball of $H(P), i(K)$ is a compact subset of $E$.

Proof :
Let us suppose first that $E$ is a Banach space (separable since $E$ is Lusin).
First we see that $i(K)$ is closed in $E$. Actually, $i$ being continuous, it is weakly continuous. Therefore $i(K)$ is weakly compact in $E$, then closed for the topology of $E$, since it is convex.

We shall now prove that $i(K)$ is relatively compact.
Identifying the dual of $H(P)$ with $E_{2}^{\prime}(P)$, the transpose of $i$ is equal to $j$, and it suffices to prove that $j: E^{\prime} \rightarrow E_{2}^{\prime}(P)$ is a compact map. This last fact may be verified as follows :

Let $\left(x_{n}^{\prime}\right)$ a sequence of elements of $B_{E^{\prime}}$, (the unit ball of $E^{\prime}$ ), converging weakly to $x^{\prime} \in B_{E^{\prime}}$, therefore converging uniformly on every compact of $E$ ( $B_{E^{\prime}}$ being an equicontinuous set).

By Fernique's theorem :

$$
\int_{E}\|x\|_{E}^{2} P(d x)<\infty
$$

applying Lebegue's dominated convergence theorem, we obtain :

$$
\left(\int_{E}\left|\left\langle x_{n}^{\prime}-x^{\prime}, x\right\rangle_{E^{\prime}, E}\right|^{2} P(d x)\right)_{n} \longrightarrow 0
$$

Therefore

$$
j\left(x_{n}^{\prime}\right) \longrightarrow j\left(x^{\prime}\right) \text { in } E_{2}^{\prime}(P)
$$

- In this case the proposition is proven. -

In the general case, let us give a sketch of the proof :
$E$ is a projective limit of a family $\left(E_{i}\right)$ of Banach spaces, (indexed by a filtering set). Let $\pi_{i}: E \rightarrow E_{i}$ (we can suppose that every $\pi_{i}$ is dense) and let $P_{i}=\pi_{i}(P)$. We know that $\pi_{i}(K)$ is the unit ball of $H\left(P_{i}\right) \subset E_{i}$.

It suffices to notice that the unit ball $K$ of $H(P)$ is compact if and only if the $\pi_{i}(K)$ are. But it is the case by the first part of the proof.

- Q.E.D. -

Proposition 4 : Let $K$ be a Hilbert space and $u: E \rightarrow K$ continuous linear. Then $u_{\mid H(P)}$ is an Hilbert-Schmidt mapping from $H(P)$ into $K$.

## Proof :

There exists a Banach $F$, such that $u$ has the following factorization :

$$
\begin{gathered}
E \xrightarrow{\pi} F \xrightarrow{\beta} K \\
u=\beta \circ \pi
\end{gathered}
$$

with $\beta$ and $\pi$ linear continuous, $\pi$ a dense map.
Let $Q=\pi(P) ; Q$ is Gaussian and $r$ sends $H(P)$ onto $H(Q)$. Then $u_{\mid H(P)}$ can be factorized as follows :

$$
\begin{gathered}
H(P) \xrightarrow{\pi^{\prime}} H(Q) \xrightarrow{\beta^{\prime}} K \\
u_{\mid H(P)}=\beta^{\prime} \circ \pi^{\prime}
\end{gathered}
$$

( $\pi^{\prime}$ is the application $\pi$ when it is considered as a mapping from $H(P)$ into $H(Q)$ and $\left.\beta^{\prime}=\beta_{\mid H(Q)}\right)$. It is sufficient to prove that $\beta^{\prime}$ is Hilbert-Schmidt and we can therefore suppose that $E$ is a separable Banach space.

We already know that $u_{\mid H}$ is compact. There exists then an orthonormal basis ( $e_{n}$ ) of $H(P)$, an orthonormal sequence $\left(f_{n}\right)$ in $K$ and a sequence $\left(\lambda_{n}\right)$ of real numbers, converging to zero such that :

$$
u_{\mid H}(h)=\sum_{n} \lambda_{n}\left(h, e_{n}\right)_{H} f_{n}, \quad \forall h \in H
$$

(we shortly write $H$ instead of $H(P)$ ).
It remains to prove that $\Sigma \lambda_{k}^{2}<\infty$.
Let for $n \in \mathbb{N}$ :

$$
X_{n}(\cdot)=\sum_{k=0}^{n} \lambda_{k} \tilde{e}_{k}(\cdot) f_{k}
$$

(with $\widetilde{e}_{k}=S^{-1}\left(e_{k}\right)$ ).
$X_{n}$ is a $K$-valued random variable defined on $(E, P)$. Clearly :

$$
\left\|X_{n}(\cdot)\right\|_{K}^{2}=\sum_{k=0}^{n} \lambda_{k}^{2}\left|\widetilde{e}_{k}(\cdot)\right|^{2}
$$

The random variables $\tilde{e}_{k}(\cdot)$ are independent and $\mathcal{N}(0,1)$.
I assert that there exists a constant $K_{1}$ (independent from $n$ and $c$ ) such that for every $c>0$ :

$$
\begin{equation*}
P\left\{x ;\left|\left\|X_{n}(x)\right\|_{K}^{2}-\sum_{k=0}^{n} \lambda_{k}^{2}\right| \geq c\right\} \leq \frac{K_{1}}{c^{2}} \sum_{k=0}^{n} \lambda_{k}^{4} \tag{A}
\end{equation*}
$$

In fact if we apply the Tchebychev's inequality to the centered random variable :

$$
\sum_{k=0}^{n}\left(\widetilde{e}_{k}^{2}(\cdot) \lambda_{k}^{2}-\lambda_{k}^{2}\right)=\sum_{k=0}^{n} \lambda_{k}^{2}\left(\widetilde{e}_{k}^{2}(\cdot)-1\right)
$$

we obtain (A) (the constant $K_{1}$ equals $\left.\mathbb{E}\left\{\left(\widetilde{e}_{k}^{2}(\cdot)-1\right)^{2}\right\}\right)$.
We deduce from (A) that for every $c>0$ :

$$
\begin{align*}
P\left\{x ;\left|\left\|X_{n}(x)\right\|_{K}^{2}-\sum_{k=0}^{n} \lambda_{k}^{2}\right|<c\left(\sum_{k=0}^{n} \lambda_{k}^{2}\right)^{\frac{1}{2}}\right\} & \geq 1-\frac{K_{1}}{c^{2}} \frac{\sum_{k=0}^{n} \lambda_{k}^{4}}{\sum_{k=0}^{n} \lambda_{k}^{2}} \\
& \geq 1-\frac{K_{1}}{c^{2}} \sup _{0 \leq k \leq n} \lambda_{k}^{2} \tag{B}
\end{align*}
$$

Now let

$$
\rho_{n, c}:=\sum_{k=0}^{n} \lambda_{k}^{2}-c\left(\sum_{k=0}^{n} \lambda_{k}^{2}\right)^{\frac{1}{2}}
$$

if $\rho_{n, c}>0$ then let :

$$
B_{n, c} \text { the open ball in } K \text { with center } 0 \text { and radius } \sqrt{\rho_{n, c}}
$$

and if it is not, then :

$$
B_{n, c}=\emptyset .
$$

It is clear that :

$$
\left\{y \in K, \quad\left|\|y\|_{K}^{2}-\sum_{k=0}^{n} \lambda_{k}^{2}\right|<c\left(\sum_{k=0}^{n} \lambda_{k}^{2}\right)^{\frac{1}{2}}\right\} \cap B_{n, c}=\emptyset
$$

From (B) we deduce :

$$
P\left\{x ;\left|\left\|X_{n}(x)\right\|_{K}^{2}-\sum_{k=0}^{n} \lambda_{k}^{2}\right|<c\left(\sum_{k=0}^{n} \lambda_{k}^{2}\right)^{\frac{1}{2}}\right\} \geq 1-\frac{K_{1}}{c^{2}} \sup _{k \in \mathbb{N}} \lambda_{k}^{2}
$$

Let us choose now $c$ such that :

$$
1-\frac{K_{1}}{c^{2}} \sup _{k \in \mathbb{N}} \lambda_{k}^{2}=\frac{1}{2}
$$

we have then :

$$
\begin{equation*}
P\left\{x \in E ; \quad X_{n}(x) \in B_{n, c}\right\} \leq \frac{1}{2}, \quad \text { for every } n \text { such that } B_{n, c} \neq \emptyset \tag{C}
\end{equation*}
$$

Now let us suppose that $u_{\mid H}$ is not Hilbert-Schmidt.
Then

$$
\sum_{k=0}^{n} \lambda_{k}^{2}-c\left(\sum_{k=0}^{n} \lambda_{k}^{2}\right)^{\frac{1}{2}} \xrightarrow[n \rightarrow \infty]{ } \infty
$$

Therefore

$$
B_{n, c} \neq \emptyset \text { if } n \text { is sufficiently large. }
$$

Moreover

$$
B_{n, c} \uparrow K,(n \longrightarrow \infty)
$$

But this last result contradicts (C).
Actually we shall see (Chapter two) that :

$$
\left(\sum_{n \leq k} \widetilde{e}_{n}(\cdot) e_{n}\right)_{k} \longrightarrow I d_{E}, \quad \text { almost surely }
$$

Therefore

$$
P\left\{x ; X_{n}(x) \in B_{n, c}\right\} \xrightarrow[n \rightarrow \infty]{ } 1
$$

There is a contradiction with (C), therefore :

$$
u_{\mid H} \text { is Hilbert-Schmidt. }
$$

Corollary 1 : Let $K$ be an Hilbert space and $v: K \rightarrow E^{\prime}$ linear and continuous when $E^{\prime}$ is equipped with the strong topology. Then $j \circ v: K \rightarrow E_{2}^{\prime}(P)$ is Hilbert-Schmidt.

## Proof :

Let us identify $K$ with its dual and let $v^{\prime}$ be the transpose of $v ; v^{\prime}$ maps the bidual $E^{\prime \prime}$ of $E$ into $K$. Let $Q$ be the image of $P$ by the canonical injection $E \rightarrow E^{\prime \prime}$. Identifying $E$ with a subset of $E^{\prime \prime}$, we know that :

$$
H(P)=H(Q) \subset E^{\prime \prime}
$$

But

$$
v_{\mid H(Q)}^{\prime}=v_{\mid H(P)}^{\prime}
$$

and

$$
v_{\mid H(Q)}^{\prime} \text { is Hilbert-Schmidt. }
$$

Therefore

$$
v_{\mid H(P)}^{\prime} \text { is Hilbert-Schmidt, and also } j \circ v
$$

- Corollary 1 is proven.-

Remark : Clearly $S \circ j \circ v$ is Hilbert-Schmidt from $K$ into $H(P)$.

Corollary 2 : Let $K_{1}$ and $K_{2}$ be two Hilbert spaces and $u_{1}: K_{1} \rightarrow E^{\prime}, u_{2}: E \rightarrow K_{2}$ linear continuous ( $E^{\prime}$ being equipped with the strong topology). Then :

$$
\begin{gathered}
u_{2} \circ i \circ S \circ u_{1}: K_{1} \rightarrow K_{2} \text { is nuclear } \\
\left(K_{1} \xrightarrow{u_{1}} E^{\prime} \xrightarrow{S} H \xrightarrow{i} E \xrightarrow{u_{2}} K_{2}\right) .
\end{gathered}
$$

Therefore we can speak of the trace of the above application.

We have already seen that if $F$ is a $P$-measurable vector subspace of $E$ such that $P(F)=1$, then $H \subset F$. We have seen too that if $\operatorname{dim} H=+\infty: P(H)=0$.

- Corollary 2 is proven.-

Now we shall give some "zero or one laws".
Proposition 5: Let $G$ be an additive subgroup of $E$, $P$-measurable. Then for every $a \in E$ we have

$$
P\{a+G\}=\left\{\begin{array}{l}
0 \\
1
\end{array}\right.
$$

(A trivial fact if $\operatorname{dim} E<\infty$ ).

## Proof :

We have supposed $E$ Lusin, $\operatorname{Supp} P=E$ and $\operatorname{dim} E=\infty!$ Now we may suppose that $E=\mathbb{R}^{\mathbb{N}}$, with the canonical Gaussian measure on it. We need several lemmas :

Lemma 1 : If $G$ is a $P$-measurable additive subgroup of $\mathbb{R}^{\mathbb{N}}$, containing $H(P)$ (in our case $H(P)=\ell^{2}$ ), then for every $a \in \mathbb{R}^{\mathbb{N}}$ :

$$
P\{a+G\}=\left\{\begin{array}{l}
0 \\
1
\end{array}\right.
$$

## Proof :

Let $\left(X_{n}\right)$ be the projection of index $n$ from $\mathbb{R}^{\mathbb{N}}$ into $\mathbb{R}$. Then the $\left(X_{n}\right)$ are independent and $\mathcal{N}(0,1)$ and moreover $\mathcal{B}_{\mathbb{R}^{\mathbb{N}}}$ is equal to $\sigma\left(X_{n}, n \in \mathbb{N}\right)$.

Now I assert that

$$
(a+G) \text { is P-independent from } \mathcal{B}_{\mathbb{R}^{\mathbb{N}}}
$$

Actually let:

$$
Y_{n}=\sum_{i \leq n} X_{i} e_{i} \quad \text { and } \quad Z_{n}(\cdot)=I d_{\mathbb{R}^{\mathbb{N}}}-Y_{n}
$$

$\left(\left(e_{i}\right)_{i \in \mathbb{N}}\right.$ is the canonical "basis" of $\left.\mathbb{R}^{\mathbb{N}}\right)$.
Since $e_{i}=S\left(X_{i}\right): Y_{n}$ takes its values in $H(P)$. Then, since $H(P) \subset G$ we have :

$$
a+G=\left\{Y_{n}+Z_{n} \in a+G\right\}=\left\{Z_{n} \in a+G\right\}
$$

Now it is clear that for each $n,\left\{Z_{n} \in a+G\right\}$ is independent from $\sigma\left(X_{i}, i \in n\right\}$, therefore $(a+G)$ is independent of $\sigma\left(X_{n}, n \in \mathbb{N}\right)$.

Our assertion is then proven. In particular $a+G$ is independent from itself.

- Lemma 1 is proven. -

Now we shall prove that the conclusion of the preceding lemma does not need the hypothesis $H \subset G$. For this we shall need the following lemma :

Lemma 2: Let $A \subset \mathbb{R}^{\mathbb{N}}$ be a $P$-measurable set such that $P(A)>0$. For every $h \in H$, there exists a real number $r_{h}>0$ with property :

$$
|r|<r_{h} \Rightarrow r h \in A-A .
$$

## Proof :

The result is true if $h=0$. So we suppose $\|h\|_{\ell^{2}}=1$. To $h$ there corresponds a random variable $\widetilde{h}$ in $E_{2}^{\prime}(P)$ as seen before.
Let

$$
Y(\cdot)=\tilde{h}(\cdot) h \quad \text { and } \quad Z=I d-Y
$$

They are two random variables with values in $\mathbb{R}^{\mathbb{N}}$, independent. Therefore :

$$
P(A)=P\{Y+Z \in A\}=\int_{\mathbb{R}^{\mathbb{N}}} P\{Y \in-z+A\} P_{Z}(d z)>0
$$

( $P_{Z}$ is the law of $Z$ ).
Therefore there exists some $z_{0} \in \mathbb{R}^{\mathbb{N}}$ such that :

$$
P\left\{Y \in-z_{0}+A\right\}>0
$$

Now let

$$
C:=\left\{r \in \mathbb{R}, \quad r h \in-z_{0}+A\right\} .
$$

Clearly

$$
\gamma_{1}(C)=P\left\{Y \in-z_{0}+A\right\}>0 .
$$

Therefore $C-C$ is a neighbourhood of zero : there exists $r_{h}>0$ such that :

$$
]-r_{h}, r_{h}[\subset C-C \subset\{r \in \mathbb{R}, \quad r h \in A-A\}
$$

- Lemma 2 is proven. -


## Now we are able to prove the proposition :

Let us suppose that $P(a+G)>0$.
Then for every $h \in H$ there exists an integer $n_{h}$ such that

$$
h \in n_{h}((a+G)-(a+G))=n_{h}(G), \quad(\text { since } G-G=G)
$$

Therefore $H \subset G$, and by Lemma 1 we conclude that

$$
P(a+G)=1
$$

- Q.E.D. -

Remark : In particular if $F$ is a $P$-measurable vector subspace, $P(F)=0$ or 1 .
We can now see from another way that if $P(F)=1$, then $H \subset F$ :
In fact, if $h \in H$, there exists $r_{h}>0$ such that

$$
r h \in F-F=F, \quad \text { if } \quad|r| \leq r_{h}
$$

But $F$ being a vector space, $r h \in F$ for every $r$.

In the next chapter we shall see that the reproducing kernel space is the intersection of the subspaces carrying $P$.

Definition : We shall call "abstract Wiener space" (or shortly "Wiener space") a triple $(E, H, P)$ where $E$ is a Lusin l.c.s., $P$ is a centered Gaussian measure on $E$ and $H=H(P)$ is the Cameron-Martin space of $P$.

## 3 - Some examples

## First let us give a trivial example :

Let $E=\mathbb{R}^{n}$ and $p \leq n$, a natural number ; let $\gamma_{p}$ the normal measure on $\mathbb{R}^{p}$ having the density $(2 \pi)^{\frac{-p}{2}} \exp \left\{-\frac{1}{2}\|x\|^{2}\right\}$. Let $P$ the image of $\gamma_{p}$ by the canonical injection $\mathbb{R}^{p} \rightarrow \mathbb{R}^{n}$. $P$ is a Gaussian measure on $\mathbb{R}^{n}$. Moreover :

$$
E_{2}^{\prime}(P)=\mathbb{R}^{p} \text { and } H(P)=\operatorname{Supp} P
$$

(Therefore $P(H(P))=1, H(P)$ is finite dimensional )
$H(P)$ is dense in $\mathbb{R}^{n}$ if and only if $p=n$.
Now let us give some less trivial examples.
Before doing this, we shall recall the notion of reproducing kernel Hilbert space of a centered real Gaussian process $\left(X_{t}\right)_{t \in T}$, indexed by an arbitrary set $T$ :
Let us denote $K$ the covariance :

$$
K(s, t)=\mathbb{E}\left\{X_{s} X_{t}\right\}
$$

Then $\mathcal{H}(T, K)$ is the subspace of $\mathbb{R}^{T}$ generated by the functions : $K(s, \cdot), s(s \in T)$.
If $u=\sum_{i=1}^{m} \alpha_{i} K\left(s_{i}, \bullet\right)$ and $v=\sum_{j=1}^{n} \beta_{j} K\left(t_{j}, \bullet\right)$ are elements of $\mathcal{H}(T, K)$ we set :

$$
(u, v)_{\mathcal{H}(T, K)}:=\sum_{i j} \alpha_{i} \beta_{j} K\left(s_{i}, t_{j}\right)
$$

It is easily proven that the right member is independent from the representations chosen of $u$ and $v$, and that $(\cdot, \cdot)_{\mathcal{H}}$ is symmetric non-negative definite (not necessarily a scalar product). $\mathcal{H}(T, K)$ is therefore a prehilbertian vector space, whose the associated Hilbert space is called the "reproducing kernel Hilbert space of the process $\left(X_{t}\right)$ " and is denoted by $H(T, K)$. Moreover, the elements in $H(T, K)$ are functions on $T$ (a non trivial fact, a priori) :

$$
H(T, K) \subset \mathbb{R}^{T}
$$

Note that "reproducing property" is satisfied, that is to say :

$$
(f(\cdot), K(s, \bullet))_{H(T, K)}=f(s), \quad \forall s \in T, \quad \forall f \in H(T, K)
$$

This relation generalizes this one : $(K(s, \bullet), K(t, \bullet))=K(s, t)$.
For all what precedes about reproducing kernel spaces of a Gaussian process, we can see Neveu [ Séminaire, Université de Montréal, Eté 1968 ], for instance.

The situation "abstract Wiener space" is an important example of what precedes. If $(E, H, P)$ is an abstract Wiener space then the process indexed by $E^{\prime}$ :

$$
x^{\prime} \leadsto\left\langle x^{\prime}, \bullet\right\rangle_{E^{\prime}, E} \text { is a Gaussian process }
$$

with covariance

$$
\left(x^{\prime}, y^{\prime}\right) \rightsquigarrow \mathbb{E}\left\{\left\langle x^{\prime}, \bullet\right\rangle\left\langle y^{\prime}, \bullet\right\rangle\right\} .
$$

Then it is easy to see that the Cameron-Martin space of $P$ is identical to the reproducing kernel space of the process.

Now we give some more concrete examples of this situation.
Let us begin by a general consideration :
Let $T$ be a compact metrisable space and let $E=\mathcal{C}(T)$ ( $E$ is a separable Banach space). Let $P$ be a Gaussian probability on $\mathcal{C}(T)$. Then the process $\left(X_{t}\right)_{t \in T}$ defined by :

$$
X_{t}(\omega)=\omega_{t} \text { if } \omega \in \mathcal{C}(T)
$$

is a Gaussian process under $P$. In this case it is easily proven that :

- for $t \in T: S\left(\delta_{t}\right)=K(t, \bullet)$, (if $K(s, t)=\mathbb{E}\left\{X_{s} X_{t}\right\}$ )
- if $\mu$ is a real measure on $T: S(\mu)(\bullet)=\int_{T} K(s, \bullet) d \mu(s)$
- if $x=S(\mu), y=S(\nu)$ are elements of $H(P)$ then

$$
(x, y)_{H(P)}=\int_{T \times T} K(s, t) d \mu(s) d \nu(t)
$$

We shall apply what precedes to the following cases.

## Example 1: Brownian motion on $[0,1]$

$P$ is the Wiener measure on $\mathcal{C}([0,1])$, carried by the space $\mathcal{C}_{0}([0,1])$ of functions null to zero, and continuous. The covariance of the process is

$$
K(s, t)=s \wedge t
$$

The Cameron-Martin space is the completion of the space of functions on $[0,1]$ from the form :

$$
\int_{[0,1]}(s \wedge \cdot) d \mu(s)
$$

where $\mu$ is a real measure on $[0,1]$.
Now if $F_{\mu}$ is the distribution function of $\mu$, by the generalised integration by parts formula, we have, for every $t$ :

$$
\begin{aligned}
\int_{[0,1]}(s \wedge t) d F_{\mu}(s) & =\left[(s \wedge t) F_{\mu}(s)\right]_{0}^{1}-\int_{0}^{1}(s \wedge t)^{\prime} F_{\mu}(s) d s \\
& =t F_{\mu}(1)-\int_{0}^{t} F_{\mu}(s) d s \\
& \left.\left.=\int_{0}^{t} \mu(] s, 1\right]\right) d s
\end{aligned}
$$

Therefore

$$
\left.\left.S(\mu)(\bullet)=\int_{0}^{\bullet} \mu(] s, 1\right]\right) d s
$$

Now the $H(P)$-norm of $S(\mu)$ is equal to :

$$
\begin{aligned}
\int_{[0,1]^{2}}(s \wedge t) d \mu(s) d \mu(t) & =\int_{[0,1]} d \mu(t)\left(\int_{[0,1]}(s \wedge t) d F_{\mu}(s)\right) \\
& \left.\left.=\int_{[0,1]}\left(\int_{0}^{t} \mu(] s, 1\right]\right) d s\right) d F_{\mu}(t) \quad \text { (by what precedes). }
\end{aligned}
$$

Integrating by parts once more, the above expression is seen equal to :

$$
\begin{aligned}
& \left.\left.\left.\left[F_{\mu}(t) \int_{0}^{t} \mu(] s, 1\right] d s\right]_{0}^{1}-\int_{0}^{1} F_{\mu}(t) \mu(] t, 1\right]\right) d t \\
& \left.\left.\left.\left.=F_{\mu}(1) \int_{0}^{1} \mu(] s, 1\right]\right) d s-\int_{0}^{1} F_{\mu}(t) \mu(] t, 1\right]\right) d t \\
& \left.\left.=\int_{0}^{1} \mu(] t, 1\right]\right)^{2} d t
\end{aligned}
$$

The function $\left.\left.\int_{0}^{\bullet} \mu(] s, 1\right]\right) d s$ is null at zero, absolutely continuous ; its derivative is Lebesguealmost everywhere equal to $\left.\left.\mu(]_{\bullet}, 1\right]\right)$.

Now let us notice that the functions $s \leadsto \mu(] s, 1])$ are dense in $L^{2}([0,1], d t)$, when $\mu$ describes the set of real measures on $[0,1]$ We can conclude that the Cameron-Martin space of the Wiener measure is the vector space of functions from $[0,1]$ into $\mathbb{R}$, absolutely continuous, vanishing at zero and whose the distributional derivative belongs to $L^{2}([0,1], d t)$, equipped with the scalar product :

$$
(f, g) \leadsto \int \dot{f}(t) \dot{g}(t) d t
$$

## Example 2 : The Brownian bridge

Now let $P$ the measure on $\mathcal{C}([0,1])$ corresponding to the Brownian bridge on $[0,1]$. The covariance of the Brownian bridge process is :

$$
(s, t) \rightsquigarrow K(s, t)=s \wedge t-s t
$$

For $\mu$, a real measure on $[0,1]$ we have for each $t$ :

$$
\left.\left.\left.\int_{[0,1]}(s \wedge t-s t) d \mu(s)=\int_{[0, t]} \mu(] s, 1\right] d s-t \int_{[0,1]} \mu(] s, 1\right]\right) d s
$$

Moreover :

$$
\left.\left.\left.\left.\int_{[0,1]^{2}}(s \wedge t-s t) d \mu(s) d \mu(t)=\int_{0}^{1}(\mu(] s, 1]\right)\right)^{2} d s-\left(\int_{0}^{1} \mu(] s, 1\right]\right) d s\right)^{2}
$$

Now let us notice that the function

$$
\left.\left.\left.\left.t \rightsquigarrow \int_{0}^{t} \mu(] s, 1\right]\right) d s-t \int_{0}^{1} \mu(] s, 1\right]\right) d s \text { is null for } t=0 \text { and } t=1
$$

It is absolutely continuous and with derivative equal almost everywhere to :

$$
\left.\left.F(t)=\mu(l t, 1])-\int_{0}^{1} \mu(] s, 1\right]\right) d s
$$

The $L^{2}(d t)$ norm of $F$ is equal to :

$$
\left.\left.\int_{[0,1]}(\mu(l t, 1])\right)^{2} d t-\left(\int_{0}^{1} \mu(l t, 1]\right) d t\right)^{2}
$$

From this we can easily see that the Cameron-Martin space of $P$ is the space of absolutely continuous functions on $[0,1]$ vanishing at zero and one, and whose derivative belongs to $L^{2}([0,1])$ equipped with scalar product :

$$
(f, g) \leadsto \int \dot{f}(t) \dot{g}(t) d t
$$

## Example 3 : The Brownian motion on $\mathbb{R}_{+}$

Let $P$ be the Gaussian measure on $E=\mathcal{C}\left(\mathbb{R}_{+}\right)$associated to it. $\mathcal{C}\left(\mathbb{R}_{+}\right)$will be equipped with the topology of uniform convergence on compacts sets : it is a separable Fréchet space (not Banach), whence a Lusin space.
$E^{\prime}$ is the space of signed Radon measures on $\mathbb{R}_{+}$with compact support.
Then as in Example $1, H(P)$ is equal to the space of functions $[0, \infty[\longrightarrow \mathbb{R}$ absolutely continuous, vanishing at zero, and whose derivatives belong to $L^{2}\left(\mathbb{R}_{+}, d t\right)$, with the natural norm

$$
\|f\|=\|\dot{f}\|_{L^{2}\left(\mathbb{R}_{+}, d t\right)}
$$

Example 4 : $E=\mathcal{S}^{\prime}\left(\mathbb{R}^{n}\right)$ (resp. $\mathcal{D}^{\prime}\left(\mathbb{R}^{n}\right)$ ) with its usual topology.
It is a Lusin space (not Fréchet). Then $E^{\prime}=\mathcal{S}\left(\mathbb{R}^{n}\right)$ (resp. $\mathcal{D}\left(\mathbb{R}^{n}\right)$ ).
By Minlos' theorem there exists a Gaussian probability $P$ on $E$ such that for every $\varphi \in E^{\prime}$ :

$$
\langle\varphi, \cdot\rangle_{E^{\prime}, E} \text { is Gaussian }
$$

with variance

$$
\|\varphi\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}
$$

The Cameron-Martin space is the space $L^{2}\left(\mathbb{R}^{n}, d t\right)$, considered as a subspace of $E^{\prime}$ by the canonical injection

$$
L^{2}\left(\mathbb{R}^{n}\right) \rightarrow \mathcal{S}^{\prime}\left(\mathbb{R}^{n}\right)
$$

Example 5: $\left(E_{i}, H_{i}, P_{i}\right)_{1 \leq i \leq n}$ is a finite family of abstract Wiener spaces.
Then the following is a Wiener space :

$$
E:=\prod_{i=1}^{n} E_{i} ; \quad H:=\bigoplus_{i=1}^{n} H_{i} ; \quad P:=\bigotimes_{i=1}^{n} P_{i}
$$

( $E$ is equipped with the product topology. It is a Lusin space).
$P$ is clearly a Radon measure, since $E$ is Lusin.
We have to prove that $P$ is Gaussian, or what is the same, that for every $x^{\prime} \in E^{\prime}$, the random variable on $E, x^{\prime}(\bullet)$ is Gaussian.
But we have (algebraically) :

$$
E^{\prime}=\prod_{i=1}^{n} E_{i}^{\prime}
$$

If $x_{i}^{\prime} \in E_{i}^{\prime}, \quad(1 \leq i \leq n)$, let us denote by $X_{i}$ the random variable :

$$
x \leadsto X_{i}(x)=x_{i}^{\prime}\left(x_{i}\right), \text { if } x=\left(x_{1}, \ldots, x_{n}\right) \in E .
$$

The $\left(X_{i}\right)(1 \leq i \leq n)$ are $P$-Gaussian and $P$ independent, and if $x^{\prime}=\left(x_{1}^{\prime}, x_{2}^{\prime}, \ldots, x_{n}^{\prime}\right)$ we have :

$$
x^{\prime}(\cdot)=\sum_{i=1}^{n} X_{i}(\cdot)
$$

This is a Gaussian random variable.
Moreover :

$$
\begin{aligned}
\left\|x^{\prime}(\cdot)\right\|_{E_{2}^{\prime}(P)}^{2} & =\sum_{i=1}^{n}\left\|X_{i}(\cdot)\right\|_{L^{2}(E, P)}^{2} \\
& =\sum_{i=1}^{n}\left\|x_{i}^{\prime}(\cdot)\right\|_{\left(E_{i}\right)_{2}^{\prime}\left(P_{i}\right)}^{2} .
\end{aligned}
$$

From this we conclude that

$$
E_{2}^{\prime}(P)=\bigoplus_{i=1}^{n}\left(E_{i}\right)_{2}^{\prime}\left(P_{i}\right)
$$

Considering each $\left(E_{i}\right)_{2}^{\prime}\left(P_{i}\right)$ as a subspace of $E_{2}^{\prime}(P)$ and each $E_{i}$ as a subspace of $E$, and apply the barycenter isomorphism $S$ of $E_{2}^{\prime}(P)$ onto $H(P)$, we can see that

$$
S \text { maps } E_{2}^{\prime}\left(P_{i}\right) \text { onto } H_{i}\left(P_{i}\right)
$$

The assertion is proven.

Example 6 : $\left(E_{n}, H_{n}, P_{n}\right)$ is a sequence of Wiener spaces.
Then :

$$
\left(\prod_{n \in \mathbb{N}} E_{n}, \bigoplus_{n \in \mathbb{N}} H_{n}, \bigotimes_{n \in \mathbb{N}} P_{n}\right) \text { is a Wiener space, too. }
$$

The proof is almost identical to the preceding one. It suffices to notice that $\left(\prod E_{n}\right)^{\prime}$ is algebraically isomorphic to the direct sum $\bigoplus_{n} E_{n}^{\prime}$ and to complete $\bigoplus_{n} H_{n}^{\prime}$ for the canonical prehilbertian structure.

Example 7 : $(E, H, P)$ is a Wiener space and $E_{0}$ is a Borel vector subspace, carrying $P$.

If $P_{0}$ denotes the induced probability by $P$ on $E_{0}$, then $\left(E_{0}, H, P_{0}\right)$ is a Wiener space. Actually, $E_{0}$, being Borelian, is a Lusin space for the induced topology. Moreover :

$$
H \subset E_{0}
$$

since $P\left(E_{0}\right)=1$.
Now each linear continuous form on $E_{0}$, is the restriction to $E_{0}$ of a linear continuous form on $E$, therefore it is Gaussian (since $P\left(E_{0}\right)=1$ ).

The announced result now follows.

Example 8 : Brownian motion with values in $\mathbb{R}^{\mathbb{N}}$ (or infinite dimensional Brownian motion)

Let $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \in \mathbb{R}^{+}}, P\right)$ a process basis and let $\left(B_{t}^{n}\right)$ a sequence of independent real Brownian motions. We shall denote

$$
\mathbb{B}_{t}:=\left(B_{t}^{n}\right)_{n} \quad t \in[0, \infty[.
$$

Each real Brownian motion determines the usual Wiener measure on $\mathcal{C}_{0}([0, \infty[, \mathbb{R})$.
Therefore $\mathbb{B}$ defines a Gaussian measure on $\left[\mathcal{C}_{0}\left([0, \infty[, \mathbb{R})]^{\mathbb{N}} \simeq \mathcal{C}_{0}\left(\left[0, \infty\left[, \mathbb{R}^{\mathbb{N}}\right)\right.\right.\right.\right.$, namely the product measure. Let $\mu_{1}$ denotes this probability ; $\mathcal{C}_{0}\left(\left[0, \infty[, \mathbb{R})^{\mathbb{N}}\right.\right.$ will be equipped with the product topology : it is a Lusin space.
$\mathcal{H}\left(\mu_{1}\right)$, the Cameron-Martin space of $\mu_{1}$ is the Hilbertian sum of the "ordinaries" CameronMartin on $\mathbb{R}$. Therefore :

$$
f \in \mathcal{H}\left(\mu_{1}\right) \Longleftrightarrow f=\left(f_{n}\right)_{n} \quad \text { and } \quad \sum_{n} \int_{0}^{\infty}\left|\dot{f_{n}}(t)\right|^{2} d t<\infty
$$

$\mathcal{H}\left(\mu_{1}\right)$ is then the subspace of $W^{1,2}\left(\left[0, \infty\left[; \ell^{2}\right)\right.\right.$ constitued by the functions null at zero. (The notation $W^{1,2}$ comes from the theory of Sobolev spaces).

Let us notice that, for every $x=\left(x_{n}\right) \in \ell^{2}$ and every $t \in[0, \infty[$ :

$$
\sum x_{n} B_{n}(t) \text { exists (almost surely) }
$$

and that

$$
t \leadsto \sum x_{n} B_{n}(t) \text { is a real Brownian motion with variance } t\|x\|_{\ell^{2}}^{2} .
$$

Now we give an important property of the infinite dimensional Brownian motion which will permit us to define the Brownian motion of an abstract Wiener space.

Lemma : Let on $\mathbb{R}^{\mathbb{N}}$ the canonical normal measure $\gamma_{1}^{\otimes \mathbb{N}}$ and let $K$ be a compact, disked ( $=$ convex, balanced) of $\mathbb{R}^{\mathbb{N}}$ such that $\gamma_{1}^{\otimes \mathbb{N}}(K)>0$ and let $F$ be the vector subspace of $\mathbb{R}^{\mathbb{N}}$ generated by $K$. Then $P$-almost surely the paths of $\mathbb{B}$ are contained in $F$ :

$$
P\left\{\omega ; \mathbb{B}_{t}(\omega) \in F \quad \forall t\right\}=1
$$

Before proving this result let us make the following remarks :
a) If $\gamma_{1}^{\otimes \mathbb{N}}(K)>0$, then $\gamma_{1}^{\otimes \mathbb{N}}(F)=1$. This results from the zero-one law. Therefore :

$$
\forall t: \quad P\left\{\omega, \mathbb{B}_{t}(\omega) \in F\right\}=1
$$

(a result weaker than the conclusion of the lemma).
b) There exists disked compacts of $\mathbb{R}^{\mathbb{N}}$ with positive measure. Actually, if $K$ is an arbitrary compact, its disked hull is compact (since $\mathbb{R}^{\mathbb{N}}$ is complete).

## Proof :

Let us suppose $t \rightsquigarrow \mathbb{B}_{t}$ continuous (this does not loss generality) and let $q_{K}$ be the gauge of $K$ :

$$
q_{K}(x)=\inf \{\lambda, \quad \lambda>0 \quad x \in \lambda K\}
$$

$q_{K}$ is convex, lower semi-continuous, with values in $[0, \infty]$.
Then $t \leadsto q_{K}^{2}\left(\mathbb{B}_{t}\right)$ is a sub-martingale with values in $[0, \infty]$, l.s.c. (but not c.à d. l.à g. in general ) : it is a composition of continuous function with l.s.c. function.

Let $T<\propto$ and $D \subset[0, T]$, denumerable, dense in $[0, T]$ and such that $T \in D$. Then :

$$
\forall t \in[0, T]: q_{K}^{2}\left(\mathbb{B}_{t}\right) \leq \liminf _{\substack{t_{n} \in D \\ t_{n} \rightarrow t}} q_{K}^{2}\left(\mathbb{B}_{t_{n}}\right) \leq \sup _{t \in D} q_{K}^{2}\left(\mathbb{B}_{t}\right)
$$

Therefore

$$
\sup _{t \leq T} q_{K}^{2}\left(\mathbb{B}_{t}\right)=\sup _{t \in D} q_{K}^{2}\left(\mathbb{B}_{t}\right)
$$

and

$$
\mathbb{E}\left\{\sup _{t \leq T} q_{K}^{2}\left(\mathbb{B}_{t}\right)\right\} \leq \mathbb{E}\left\{\sup _{t \in D} q_{K}^{2}\left(\mathbb{B}_{t}\right)\right\} \leq 4 \mathbb{E}\left\{q_{K}^{2}\left(\mathbb{B}_{T}\right)\right\}
$$

(by Doob's inequality). Now :

$$
4 \mathbb{E}\left\{q_{K}^{2}\left(\mathbb{B}_{T}\right)\right\}=4 T \mathbb{E}\left\{q_{K}^{2}\left(\mathbb{B}_{1}\right)\right\} \leq 4 T \int_{\mathbb{R}^{\mathbb{N}}} q_{K}^{2} d \gamma_{1}^{\otimes \mathbb{N}}
$$

But, by Fernique's theorem we have :

$$
\int_{\mathbb{R}^{\mathbb{N}}} q_{K}^{2} d_{r_{1}}^{8 \mathbb{N}}=\int_{F} q_{K}^{2} d^{F}\left(\gamma_{1}^{\otimes N}\right)<\infty
$$

where $d^{F}\left(\gamma_{1}^{\otimes \mathbb{N}}\right)$ denotes the induced measure on $F$ (which carries $\gamma_{1}^{\otimes \mathbb{N}}$ ).
Finally, for almost every $\omega,\left(\mathbb{B}_{t}\right)$ stays in an homothetic of $K$ when $t \in[0, T] . T$ being arbitrary :

- the lemma is proven. -


## Example 9 : Brownian motion in an abstract Wiener space

Let $(E, H, P)$ a Wiener space. We shall suppose fulfilled the following condition :
"There exists a compact disk $K$ of $E$ such that $P(K)>0$ ".
Now let $\pi: E \rightarrow \mathbb{R}^{\mathbb{N}}$ a linear continuous injection of $E$ into $\mathbb{R}^{\mathbb{N}}$ mapping $P$ on $\gamma_{1}^{\otimes \mathbb{N}}$, and let $L=\pi(K)$, where $K$ is as above.

The vector subspace generated by $L$ (v.s.(L)) carries $\gamma_{1}^{\otimes \mathbb{N}}$ as we just saw and it is clearly contained in $\pi(E)$. Moreover, we know that, almost surely the paths of the canonical Brownian motion with values in $\mathbb{R}^{\mathbb{N}}$ are contained in v.s.( L$)$. Then $\pi^{-1}\left(\mathbb{B}_{t}\right)=W_{t}$ has sense, and defines a $E$-valued process.

Now we can see easily that if $f \in E^{\prime}$, there exists an $x \in \ell^{2}$ associated to it such that

$$
\|f\|_{E_{2}^{\prime}(P)}=\|x\|_{\ell^{2}}
$$

and that :

$$
f \circ W_{t} \text { is a Brownian motion (real) }
$$

with variance

$$
t\|f\|_{E_{2}^{\prime}(P)}^{2}
$$

This result remains true if $f \in E_{2}^{\prime}(P)$.
Therefore we proved the following theorem :

THEOREM 3: Let $E$ be a Lusin l.c.s. and let $P$ be a Gaussian probability on $E$ satisfying the above condition. There exists a process $\left(W_{t}\right)_{t \in \mathbb{R}_{+}}$with values in $E$ such that, for every $f \in E_{2}^{\prime}(P),\left(f \circ W_{t}\right)_{t \in \mathbb{R}_{+}}$is a real Brownian motion with variances $t\|f\|_{E_{2}^{\prime}(P)}$.
$\left(W_{t}\right)$ is called the "canonical Brownian motion" associated to $P$, or to $(E, P)$.
Remark : The condition of the theorem is automatically fulfilled when $E$ is Banach, or complete, or more generally quasi-complete.

Now let us end by a characterization of a Gaussian probability which will be used in several places:

A Borelian probability $P$ on a Lusin l.c.s. $E$ is Gaussian if and only if it satisfies the following property :
"If $X$ and $Y$ are E-valued random variables (defined on a same probability space) independent and having the common law $P$, then for every $\theta \in \mathbb{R}$, the $E$-valued random variable $X \cos \theta+Y \sin \theta$ has the law $P$ ".

This property is immediately deduced from the corresponding characteristics of a real centered Gaussian law.

This property was used by Fernique to define a Gaussian probability on a general topological vector space (not necessarily locally convex).

## CHAPTER TWO

## Measurable linear operators

## 1 - Generalities

Let us recall some fundamental notions.
Let $X$ be a Hausdorff completely regular topological space and let $\mu$ be a Borelian probability on $X$. Let $Y$ be another Hausdorff topological space. Then we have three kinds of measurability of a linear application $f: X \rightarrow Y$ :
a) the "Borel-measurability" (does not depend of $\mu$ ),
i .e. the measurability with respect to the Borelian $\sigma$-fields $\mathcal{B}_{X}$ and $\mathcal{B}_{Y}$,
b) the "Lusin-measurability":
$f: X \rightarrow Y$ is said $\mu$-Lusin measurable if it satisfies the following condition :
for every $\varepsilon>0$ there exits a compact $K_{\varepsilon} \subset X$ such that :

$$
\mu\left(X \backslash K_{\varepsilon}\right)<\varepsilon \quad \text { and } \quad f_{\mid K_{\varepsilon}} \text { is continuous }
$$

c) the " $\mu$-measurability", or "Lebesgue-measurability":
$f: X \rightarrow Y$ is said $\mu$-measurable if it is measurable with respect to $\mathcal{B}_{X}^{\mu}$ (the $\mu$-completed $\sigma$-field of $\mathcal{B}_{X}$ ) and $\mathcal{B}_{Y}$.

Clearly :

$$
\text { a) } \Longrightarrow c) \quad \text { and } \quad \text { b) } \Longrightarrow c)
$$

Moreover :

$$
\text { if } \mu \text { is Radon, then } \mathrm{c}) \Longrightarrow \mathrm{b} \text { ). }
$$

Actually this fact is well known if $X$ is locally compact.
In the general case we can do the following reasoning :
Let $L_{\varepsilon}$ be a compact such that :

$$
\mu\left(L_{\varepsilon}\right) \leq \varepsilon / 2
$$

By the Egoroff's theorem there exists a compact $K_{\varepsilon} \subset L_{\varepsilon}$ such that :

$$
\mu\left(L_{\varepsilon} \backslash K_{\varepsilon}\right)<\varepsilon / 2 \quad \text { and } \quad f_{\mid K_{\varepsilon}}=\left(f_{\mid L_{\varepsilon}}\right)_{\mid K_{\varepsilon}} \text { be continuous. }
$$

Then

$$
\mu\left(X \backslash K_{\varepsilon}\right)<\varepsilon
$$

and the result is proven.

Let us suppose now that $Y$ is a Hausdorff l.c.s.; we can then define three kinds of scalar measurability :
$f: X \rightarrow Y$ is said scalarly measurable (resp. $\mu$-measurable, resp. Lusin-measurable) if : for every linear continuous form $y^{\prime}: y^{\prime} \circ f$ is measurable (resp. $\mu$-measurable, resp. Lusin-measurable).

Lemma 1 : Let $(X, \mu)$ as above, with $\mu$ Radon and let $Y$ be a Lusin l.c.s. (Hausdorff). If $f: X \rightarrow Y$ is scarlarly $\mu$-measurable, then it is $\mu$-measurable.
(Therefore we have the equivalence : $\mu$-measurability $\Longleftrightarrow$ scalar $\mu$-measurability).
Proof :
Let $j$ be a linear continuous injection of $Y$ into $\mathbb{R}^{\mathbb{N}}$, then :

$$
j(Y) \text { is a Borel set of } \mathbb{R}^{\mathbb{N}}
$$

and :

$$
j: Y \rightarrow j(Y) \text { is bi-measurable }
$$

(with respect to the Borel $\sigma$-fields).
Let $g=j \circ f, g: X \rightarrow \mathbb{R}^{\mathbb{N}}$, then $g$ is $\mu$-measurable by the definition of the topology of $\mathbb{R}^{\mathbb{N}}$. Therefore :

$$
j^{-1} \circ j \circ f \quad \text { is } \mu \text {-measurable }
$$

since $j^{-1}$ is Borelian.

- Lemma 1 is proven.-

We have a similar result for Lusin-measurability.
In what follows $X$ and $Y$ will be Hausdorff locally convex spaces ("l.c.s.") and $f$ will be linear (often denoted $T$ ).

Lemma 2: Let $E$ be a Hausdorff l.c.s., $\mu$ be a Radon probability on $E$ and $T: E \rightarrow F$ linear Lusin-measurable (where $F$ is a topological vector space). Then there exists a Borel subspace $E_{0}$ of $E$, carrying $\mu$ and such that $T_{\mid E_{0}}$ is Borelian.

## Proof :

Let $\varepsilon>0$ and $K_{\varepsilon}$ be a compact (not necessarily convex) such that :

$$
\mu\left(K_{\varepsilon}\right) \geq 1-\varepsilon \quad \text { and } \quad T_{\mid K_{\varepsilon}} \text { be continuous, }
$$

and let $E^{\varepsilon}$ be the vector subspace generated by $K_{\varepsilon}$.
I assert that $E^{\varepsilon}$ is Borelian and $T_{\mid E^{\varepsilon}}$ is a Borelian map. In fact,let us fix $n \in \mathbb{N}$ and let :

$$
K_{\varepsilon}^{n}:=\left\{x \in E, \quad x=\sum_{i=0}^{i=n} \lambda_{i} x_{i}, \quad\left|\lambda_{i}\right| \leq 1, \quad x_{i} \in K_{\varepsilon}, \quad \forall i\right\}
$$

First $K_{\varepsilon}^{n}$ is compact because : if $\left(x^{\alpha}\right)_{\alpha}$ is an ultrafilter on $K_{\varepsilon}^{n}$, it converges to $x \in K_{\varepsilon}^{n}$. Actually :

$$
x^{\alpha}=\sum_{i=0}^{i=n} \lambda_{i}^{\alpha} x_{i}^{\alpha}
$$

then :

$$
\lambda_{i}^{\alpha} \underset{\alpha}{\longrightarrow} \lambda_{i}, \quad(0 \leq i \leq n)
$$

since $[-1,+1]$ is a compact of $\mathbb{R}$ and

$$
x_{i}^{\alpha} \underset{\alpha}{\longrightarrow} x_{i}
$$

since $K_{\varepsilon}$ is compact.
Moreover : $T_{\mid K_{\epsilon}^{n}}$ is continuous, therefore :

$$
T\left(x^{\alpha}\right)=\sum_{i=0}^{i=n} \lambda_{i}^{\alpha} T\left(x_{i}^{\alpha}\right) \longrightarrow \sum_{i=0}^{i=n} \lambda_{i} T\left(x_{i}\right)
$$

Let us notice that $K_{\varepsilon}^{n}$ is not necessarily convex, but it is balanced . Then :

$$
T_{\bigcup_{m} m K_{\varepsilon}^{n}} \quad \text { is Borelian }
$$

and :

$$
T_{\bigcup_{n} \bigcup_{m}\left(m K_{e}^{n}\right)} \quad \text { is Borelian too. }
$$

We can write

$$
E^{\varepsilon}=\bigcup_{n} \bigcup_{m} m K_{\varepsilon}^{n} ;
$$

then $E_{\varepsilon}$ is Borelian and it is a vector subspace.
Now let $\varepsilon=1 / p,\left(p \in \mathbb{N}^{*}\right)$. We may choose the sequence $\left(K_{1 / p}\right)_{p}$ increasing. Then :

$$
E_{0}=\bigcup_{p \in \mathbb{N}^{*}} E^{1 / p}
$$

satisfies the condition of the lemma.

- Lemma 2 is proven. -

Remark 1 : The proof of Lemma 2 can be simplified if we know that $\mu$ is almost carried by compact convex sets. In this case :

$$
\bigcup_{m} m K_{\varepsilon} \quad \text { is a vector space }
$$

if $K_{\varepsilon}$ is disked.

Remark 2: If $E$ and $F$ are l.c.s. and $\mu$ a Radon-measure on $E$, the Lusin-measurable functions from $E$ into $F$ form a vector space (it is obvious by the definition). If $E$ and $F$ are Lusin, the Borel functions from $E$ into $F$ form a vector space. It is obvious since :

$$
\mathcal{B}_{E} \otimes \mathcal{B}_{F}=\mathcal{B}_{E \times F}
$$

and :

$$
\mathcal{B}_{\mathbb{R} \times E}=\mathcal{B}_{\mathbb{R}} \otimes \mathcal{B}_{E}
$$

In the sequel, unless the contrary is specified, we shall assume that the vector spaces we consider are Lusin.

Lemma 3 : Let $\left(E_{1}, \mu_{1}\right)$ and $\left(E_{2}, \mu_{2}\right)$ be two (Lusin) l.c.s. with (Radon) probabilities. Let $F$ be another Lusin l.c.s. and $T: E_{1} \times E_{2} \rightarrow F$ linear, $\mu_{1} \otimes \mu_{2}$-measurable. Then for every $x_{2} \in E_{2}$, the map $T\left(\cdot, x_{2}\right)$ is $\mu_{1}$-measurable.

## Proof :

By Fubini's theorem :
for $\mu_{2}$-almost every $x_{2} \in E_{2}: \quad T\left(\cdot, x_{2}\right)$ is $\mu_{1}$-measurable,
there exists $x_{2}^{0} \in E_{2}$ such that :

$$
T\left(\cdot, x_{2}^{0}\right) \text { is } \mu_{1} \text {-measurable. }
$$

By the remark preceding the lemma :

$$
T(\cdot, 0)=\frac{1}{2}\left[T\left(\cdot, x_{2}^{0}\right)+T\left(\cdot,-x_{2}^{0}\right)\right] \quad \text { is } \mu_{1} \text {-measurable. }
$$

Now for each $x_{2} \in E_{2}$ :

$$
T\left(\bullet, x_{2}\right)=T(\cdot, 0)+T\left(0, x_{2}\right)
$$

- Lemma 3 is proven. -

Remark : Let $E$ be a l.c.s. with the Borel probability $\mu$ on it and let $E_{0}$ be a Borel subspace of $E$, carrying $\mu$. Let $F$ another l.c.s. and $T_{0}: E_{0} \rightarrow F$ linear, Borelian. Every linear extension of $T_{0}$ from $E_{0}$ to $E$ (easily obtained by considering an Hamel basis of $E$ which completes an Hamel basis of $E_{0}$ ) is clearly $\mu$-measurable, but not necessarily Borelian.

## This fact should be kept in mind.

## 2 - Measurable linear mappings from a Wiener space

In the sequel, $E$ will denote a Lusin l.c.s. (and Hausdorff), $\mu$ a Gaussian centered probability on $E$ whose reproducing kernel Hilbert-space will be denoted by $H$ or $H(\mu)$. We shall suppose that $\operatorname{Supp}(\mu)=E$, then $H$ is dense in $E$. Otherwise there exists $x_{0}^{\prime}$ in $E^{\prime}$ such that

$$
x_{0}^{\prime} \neq 0, x_{0}^{\prime}(h)=0 \text { for each } h \text { in } H .
$$

Therefore :

$$
0=x_{0}^{\prime}\left(S x_{0}^{\prime}\right)=\int_{E}\left|x_{0}^{\prime}(x)\right|^{2} d \mu(x)
$$

Then :

$$
\mu\left\{x ; x_{0}^{\prime}(x)=0\right\}=1
$$

which contradicts Supp $(\mu)=E$.
The results that follow are true, but trivial, if:

$$
E=H=\mathbb{R}^{n}
$$

so we shall exclude this case in our proofs and we shall suppose : $\operatorname{dim} H=+\infty$.
Let $F$ be another Lusin l.c.s., let $T: E \rightarrow F$, and $T^{\prime}: E \rightarrow F$ be two $\mu$-measurable linear maps. Let us suppose that : $T=T^{\prime} \mu$-almost surely. Then :

$$
\left\{x, \quad T(x)=T^{\prime}(x)\right\}
$$

is a subspace carrying $\mu$ and :

$$
H \subset\left\{T=T^{\prime}\right\}
$$

In other words:

Two linear measurable maps, equal almost surely, coincide on $H$.
(The converse is true, see below).
Then we can speak about the value at a point of $H$, of a $\mu$-equivalence class of linear $\mu$-measurable maps.

On the other hand, by using the caracterization of a Gaussian measure given at the end of chapter one : $T(\mu)$ is again a Gaussian measure if $T$ is $\mu$-measurable and linear.

Let us give a consequence from Lemma 2 .

THEOREM 1 : Let $F$ be a Banach and $T: E \rightarrow F$ linear $\mu$-measurable. Then we have : $\int_{E}\|T(x)\|^{2} \mu(d x)<\infty$.

Proof :
By Lemma 2, there exists a Borelian subspace, $E_{0}$, carrying $\mu$ such that $T_{\mid E_{0}}$ is Borelian. Then $E_{0}$ is a Lusin-space for the induced topology and

$$
\mathcal{B}_{E_{0} \times E_{0}}=\mathcal{B}_{E_{0}} \otimes \mathcal{B}_{E_{0}}
$$

Now let us consider the Borelian semi-norm on $E_{0}$ :

$$
x \leadsto\|T x\|_{F}
$$

By the Fernique's theorem (see chapter one) we have :

$$
\int_{E_{0}}\|T x\|_{F}^{2} \mu_{0}(d x)<\infty
$$

where $\mu_{0}=\mu_{\mid E_{0}}$ and therefore :

$$
\int_{E}\|T x\|_{F}^{2} \mu(d x)<\infty
$$

- Q.E.D.-

First we examine the case : $F=\mathbb{R}$, i.e. we shall consider the $\mu$-measurable linear forms.

Let us notice that we have met such linear forms : the linear forms defined by an element of $E_{2}^{\prime}(\mu)$. In fact such an element is an almost sure limit of a sequence of linear continuous forms, whose the convergence set is a Borelian vector space carrying $\mu$. Then we extend it outside of this vector subspace with preservation of the linearity, and we obtain a linear $\mu$-measurable form .

Now we shall see that there are not other $\mu$-measurable linear forms. Actually let $f: E \rightarrow \mathbb{R}$ linear $\mu$-measurable, therefore defining an element of $L^{2}(E, \mu)$ and let us denote by $\widehat{v}$ its orthogonal projection on $E_{2}^{\prime}(\mu)$.

## I assert that :

$$
f=\widehat{v}
$$

because : $f-\hat{v}$ is orthogonal to each $\widehat{u} \in E_{2}^{\prime}(\mu)$.
Actually, let $\left(e_{n}^{\prime}\right)_{n}$ be a sequence of elements in $E^{\prime}$, which constitues an orthonormal basis of $E_{2}^{\prime}(\mu)$, the family of random variables on E :

$$
\left\{\left\langle e_{n}^{\prime}, \cdot>, n \in \mathbb{N}\right\} \bigcup\{f-\widehat{v}\}\right. \text { is Gaussian }
$$

(it defines a $\mu$-measurable linear mapping from $E$ to $\mathbb{R}^{\mathbb{N}} \otimes \mathbb{R}$ ).
Since $f-\widehat{v}$ is orthogonal to the $\left\langle e_{n}^{\prime}, \bullet\right\rangle,(n \in \mathbb{N})$, it is independent from the $\left\langle e_{n}^{\prime}, \bullet\right\rangle$ and therefore it is independent from $\mathcal{B}_{E}=\sigma\left\{\left\langle e_{n}^{\prime}, \cdot\right\rangle, \quad n \in \mathbb{N}\right\}$.
Therefore :

$$
f-\widehat{v} \text { is a constant (almost surely) }
$$

and this constant is null, since $f-\widehat{v}$ is centered.
The announced result is established.

From what we saw in the chapter one, we can see that the restriction of $f$ to $H$ is continuous. This last fact can be seen otherwise as follows :
if $E_{0}$ is a Borel subspace of $E$ carrying $\mu$ and such that $f_{\mid E_{0}}$ be Borelian then, since :

$$
H \subset E_{0} \text { with continuous injection }
$$

( $E_{0}$ being equipped with the induced topology, $H$ with the Hilbertian topology) we have :

$$
f_{\mid H} \text { Borelian and therefore continuous }
$$

(all Borelian linear forms on a Banach space being continuous).
Summing up that precedes, we have obtained the following theorem :

THEOREM 2 : Let $(E, H, \mu)$ be an abstract Wiener space. Then there is an identity between :
(i) : the $\mu$-equivalence classes of linear $\mu$-measurable forms,
(ii) : the almost sure limits of sequences of linear continuous forms.

Moreover all linear continuous forms on $H$ (with its topology) admit a linear $\mu$ measurable extension to $E$ and two such extensions belong to the same class. If $g$ is a linear continuous form, we shall call " $\mu$-essential extension of $g$ " whatever extension.

Let us notice the equality:

$$
\|g\|_{H^{\prime}}^{2} \equiv \int_{E}|f(x)|^{2} \mu(d x)
$$

if $f$ extends $g$.
Remark 1 : We have no more the $\mu$-unicity of the extension if we do not suppose the $\mu$-measurability of extensions, as we can see :
let $f$ be a linear form on $E$, not $\mu$-almost surely equal to zero and vanishing on $H$ (we can build such a linear form by mean of a suitable Hamel basis).

Remark 2 : The $\mu$-measurable linear forms, obviously separate the elements of $E$.
Now let us give a criterion for an element of $E$ to belong to $H$.
Proposition 1 : Let $a \in E$. The following properties are equivalent :
(i) $: a \in H$
(ii) : there exists $C \in\left[0, \infty\left[\right.\right.$ such that $:|f(a)| \leq C\|f\|_{L^{2}(E, \mu)}$ for each $\mu$-measurable linear form $f$.

## Proof :

$(i) \Rightarrow(i i)$ : it is trivial, we have $C=\|a\|_{H}$,
(ii) $\Rightarrow(i)$ : let us notice that if $f=g$ almost surely ( $f$ and $g$ linear $\mu$-measurable ) then, by condition (ii),

$$
f(a)=g(a) .
$$

Therefore we can define a linear mapping $E_{2}^{\prime}(\mu) \rightarrow \mathbb{R}$ by:

$$
f \leadsto f(a) .
$$

This linear form is continuous (with norm $\leq C$ ). So there exists $u \in H$ such that :

$$
\text { for every } f \text { in } E_{2}^{\prime}(\mu): \quad f(a)=f(u)=(S f, u)_{H}
$$

Therefore

$$
\text { for every } f \in E_{2}^{\prime}(\mu): \quad f(u-a)=0
$$

And $u=a$ since the $f$ are separating.

Corollary : Let $a \in E$, the following properties are equivalent :
(i) $: a \in H$
(ii) : for every sequence $\left(f_{n}\right)$ of $\mu$-measurable linear forms converging almost surely to zero, we have $\operatorname{bigl}\left(f_{n}(a)\right) \longrightarrow 0$.

Proof :
$(i) \Rightarrow(i i):$ if $f_{n} \longrightarrow 0$ almost surely then :

$$
f_{n} \longrightarrow 0 \text { in } L^{2}(E, \mu)
$$

(since the $f_{n}(\cdot)$ are Gaussian random variables).
Therefore

$$
f_{n}(a) \longrightarrow 0 \quad \text { by Proposition } 1
$$

$(i i) \Rightarrow(i)$ : suppose (ii) realized and suppose that $a \notin H$. Then there exists a sequence $\left(f_{n}\right)_{n}$ of $\mu$-measurable linear forms converging to zero in $L^{2}$ such that $\left(f_{n}(a)\right)$ does not converge to zero.
So there exists a subsequence $\left(f_{n_{k}}\right)$ such that :

$$
f_{n_{k}} \xrightarrow[k]{\longrightarrow} \text { in } L^{2}
$$

and

$$
\left(f_{n_{k}}(a)\right)_{k} \text { converges in } \overline{\mathbb{R}} \text { to } \alpha \neq 0
$$

From this we deduce the existence of a subsequence of $\left(f_{n_{k}}\right)_{k}$ denoting by $\left(g_{\ell}\right)_{\ell}$ converging almost surely to zero, but such that :

$$
\left(g_{\ell}(a)\right)_{\ell} \text { does not converge in } \mathbb{R} \text { to zero, }
$$

- The Corollary is proven.-

Now let us fix some notations :
If $\left(f_{n}\right)_{n}$ is a sequence of linear $\mu$-measurable forms converging almost surely to zero, we shall denote the set of convergence of $\left(f_{n}\right)_{n}$ by: $C\left\{\left(f_{n}\right)_{n}\right\}$. This is a vector subspace carrying $\mu$. So :

$$
H \subset C\left\{\left(f_{n}\right)_{n}\right\}
$$

Proposition 2 : $H$ is equal to the intersection of the $C\left\{\left(f_{n}\right)_{n}\right\}$ where $\left(f_{n}\right)_{n}$ is a sequence of elements in $E_{2}^{\prime}(\mu)$ converging almost surely to zero.

## Proof :

Let us denote by $D$ the intersection of the $C\left\{\left(f_{n}\right)_{n}\right\}$.
Clearly :

$$
H \subset D
$$

Let us suppose that there exists $a \in D \backslash H$. Since $a \notin H$, by the preceding Corollary there exists $\left(f_{n}\right)_{n} \longrightarrow 0$ almost surely such that :

$$
\left(f_{n}(a)\right)_{n} \text { does not converge to zero. }
$$

There is a contradiction.

As an immediate consequence we have the following theorem :
THEOREM 3 : $H$ is the intersection of the Borel vector subspaces carrying $\mu$.
$\diamond$ Now let us come back to the general case of a $\mu$-measurable linear mapping between two Lusin spaces.

Let $T: E \rightarrow F$ linear and $\mu$-measurable. As when $T$ is continuous we have the relation

$$
\text { for every } f \in L^{2}(E, \mu): \quad T\left(\int_{E} f(x) x \mu(d x)\right)=\int_{E} f(x) T(x) \mu(d x)
$$

where the right member is a "weak integral":

$$
\forall g \in F^{\prime}:<g, \int_{E} f(x) T(x) \mu(d x)>:=\int_{E} f(x) g \circ T(x) \mu(d x)
$$

(Let us notice that the last integral has a sense since : $g \circ T$ is a $\mu$-measurable linear functional on E ). In fact :

$$
\begin{aligned}
<g, T\left(\int_{E} f(x) x \mu(d x)\right)>_{E^{\prime}, E} & =g \circ T\left(\int_{E} f(x) x \mu(d x)\right) \\
& \left.=\int_{E} g \circ T(x) f(x) \mu(d x) \quad \text { (since } \quad \mathrm{g} \circ \mathrm{~T} \in \mathrm{E}_{2}^{\prime}(\mu)\right) \\
& =<g, \int_{E} f(x) T(x) \mu(d x)>\quad \text { ( } \mathrm{g} \text { being continuous ). }
\end{aligned}
$$

Proposition 3 : Let $\nu=T(\mu)$ and let $K$ be the Cameron-Martin space of $\nu$. Then $T(H)=K$. Moreover, $T$ induces a linear continuous function from $H$ into $K$.
(This result was already proven in chapter one, when $T$ was continuous).

## Proof :

$S_{\nu}: g \leadsto \int_{F} g(y) y d \nu$ is an isometry from $F_{2}^{\prime}(\nu)$ onto $K$.
Let $G:=\left\{\int_{E} g \circ T(x) x \mu(d x) ; g \in F_{2}^{\prime}(\nu)\right\}$. Let $g \in F_{2}^{\prime}(\nu)$, so $g \circ T \in E_{2}^{\prime}(\mu)$. By the remark preceding this proposition :

$$
\begin{aligned}
T\left(\int_{E} g \circ T(x) x \mu(d x)\right) & =\int_{E} g \circ T(x) T(x) \mu(d x) \\
& =\int_{F} g(y) y \nu(d y)
\end{aligned}
$$

So

$$
K=T(G) \subset T(H)
$$

Now we prove that $T(G)=T(H)$ :
Let $\widehat{u} \in E_{2}^{\prime}(\mu)$, orthogonal to the $g \circ T$, with $g \in F_{2}^{\prime}(\nu)$ and let $u$ be the barycenter of $\widehat{u}$. I claim that :

$$
T(u)=0
$$

In fact for $g \in F^{\prime}$ :

$$
\begin{aligned}
<g, T(u)>_{F^{\prime}, F} & =g \circ T\left(\int_{E} \widehat{u}(x) x \mu(d x)\right) \\
& =\int_{E} \widehat{u}(x) g \circ T(x) \mu(d x)=0
\end{aligned}
$$

This equality being true for every $g \in F^{\prime}$, then $T(u)=0$. So :

$$
T(G)=T(H)
$$

Now the continuity of $T_{\mid H}$ is obvious : $T_{\mid H}$ is the transpose of the isometry: $g \rightsquigarrow g \circ T$ from $F_{2}^{\prime}(\nu)$ into $E_{2}^{\prime}(\mu)$ as we can see :
let $\widehat{v} \in F_{2}^{\prime}(\nu)$ with barycenter $v$, then

$$
\begin{aligned}
\left(T\left(\int_{E} g \circ T(x) x \mu(d x)\right) ; v\right)_{K} & =\left(\int_{E} g \circ T(x) T(x) \mu(d x) ; v\right)_{K} \\
& =\left(\int_{F} g(y) y \nu(d y) ; v\right)_{K} \\
& =\int_{F} g(y) \widehat{v}(y) \nu(d y) \\
& =\int_{E} g \circ T(x) \widehat{v} \circ T(x) \mu(d x) \\
& =(g \circ T ; \widehat{v} \circ T)_{E_{2}^{\prime}(\mu)}
\end{aligned}
$$

- Q.E.D. -

We have seen that if $T_{1}=T_{2} \mu$-almost surely, then their restrictions to $H$ coincide. The converse is true.

Proposition 4 : Let $T_{1}$ and $T_{2}$ be two $\mu$-measurable linear maps from $E$ to $F$ such that $T_{1 \mid H}=T_{2 \mid H}$, then $T_{1}=T_{2} \mu$-amost surely.

## Proof :

The result is true if $F=\mathbb{R}$ (linear measurable functionals) as we have seen above. Therefore the proposition remains true if $F=\mathbb{R}^{\mathbb{N}}$.

For the general case, let us consider a continuous linear injection $j$ from $F$ into $\mathbb{R}^{\mathbb{N}}$ and the functions $j \circ T_{i}, \quad(i=1,2)$.

## 3 - Comparison of the Cameron-Martin spaces.

THEOREM 4 : Let $\left(E, H_{1}, \mu_{1}\right)$ and $\left(E, H_{2}, \mu_{2}\right)$ be two Wiener spaces (with the same $E)$; we suppose that every subspace carrying $\mu_{2}$ carries $\mu_{1}$, then:
(i) : every linear form $\mu_{2}$-measurable is $\mu_{1}$-measurable and we can determine a linear continuous function from $E_{2}^{\prime}\left(\mu_{2}\right)$ into $E_{2}^{\prime}\left(\mu_{1}\right)$
(ii) : $H_{1} \subset H_{2}$ with continuous injection.

Proof :
(i) : Let $f: E \rightarrow \mathbb{R}$ linear and $\mu_{2}$-measurable. Then there exists a Borelian subspace $F$ of $E$ such that :

$$
f_{\mid F} \text { is Borelian and } \mu_{2}(F)=1
$$

Since $\mu_{1}(F)=1$, then

$$
f \text { is } \mu_{1} \text {-measurable. }
$$

Now let us suppose that $f_{1}=f_{2}, \quad \mu_{2}$-almost surely ( $f_{1}$ and $f_{2}$ linear, $\mu_{2}$-measurable). Then $\left\{f_{1}=f_{2}\right\}$ is a subspace carrying $\mu_{2}$ and therefore $\mu_{1}$. This means

$$
f_{1}=f_{2}, \quad \mu_{1} \text {-almost surely. }
$$

So we have defined a canonical injection

$$
E_{2}^{\prime}\left(\mu_{2}\right) \rightarrow E_{2}^{\prime}\left(\mu_{1}\right)
$$

The continuity of this injection follows from the closed-graph theorem :
if $f_{n} \longrightarrow f$ in $E_{2}^{\prime}\left(\mu_{2}\right)$ and $f_{n} \longrightarrow g$ in $E_{2}^{\prime}\left(\mu_{1}\right)$ then there exists a subsequence $\left(f_{n_{k}}\right)$ converging $\mu_{2}$-almost surely to $f$.
Therefore

$$
g=f, \mu_{1} \text {-almost surely }
$$

There exists $C \in[0, \infty[$ such that :

$$
\int(f(x))^{2} \mu_{1}(d x) \leq C \int(f(x))^{2} \mu_{2}(d x), \quad \forall f \in E_{2}^{\prime}\left(\mu_{2}\right)
$$

and (i) is proven.
(ii) : The inclusion $H_{1} \subset H_{2}$ follows immediately from Theorem 3 and the continuity of this inclusion follows from the closed-graph theorem.

Remark : If $\mu_{1}$ and $\mu_{2}$ satisfy the condition of the Theorem 4, for every $F$ Lusin l.c.s., we have :

$$
T: E \rightarrow F \text {, linear, } \mu_{2} \text {-measurable } \Rightarrow T \mu_{1} \text {-measurable }
$$

(by Lemma 1).

Now we shall be given with two Wiener spaces $\left(E_{i}, H_{i}, \mu_{i}\right),(i=1,2)$ and $T: E_{1} \rightarrow E_{2}$ linear and $\mu_{1}$-measurable. Then we can easily prove the following theorem :

THEOREM 5 : Let us suppose that every vector subspace of $E_{2}$ carrying $\mu_{2}$ carries $T\left(\mu_{1}\right)$. Then $T\left(H_{1}\right) \subset H_{2}$ and $T_{\mid H_{1}}: H_{1} \rightarrow H_{2}$ is continuous.

Proof :
Let $\nu_{1}=T\left(\mu_{1}\right)$. We have noticed that $\nu_{1}$ is Gaussian. Let $K_{1}$ be the CameronMartin space of $\nu_{1}$. We know that $T$ maps continuously $H_{1}$ into $K_{1}$ (with their Hilbertian topology).

By the hypothesis every subspace carrying $\mu_{2}$ carries $\nu_{1}$, therefore $K_{1} \subset H_{2}$ with continuous injection. So

$$
T: H_{1} \rightarrow H_{2} \text { is continuous. }
$$

Corollary 1 : Under the hypothesis of the Theorem 5, the restriction of $T$ to $H_{1}$ is continuous from $H_{1}$ into $E_{2}$.

Next we shall be concerned with the inverse problem : let $T: H_{1} \rightarrow H_{2}$ linear and continuous, does there exist an extension $\widetilde{T}: E_{1} \rightarrow E_{2}$ linear and $\mu_{1}$-measurable? The answer is affirmative.
( Let us notice that if such an extension exists, it is unique up-to a $\mu_{1}$-equivalence by Proposition 4 ).

THEOREM 6 : Let $T: H_{1} \rightarrow H_{2}$ linear and continuous. Then there exists an extension $\widetilde{T}: E_{1} \rightarrow E_{2}$ of $T$, linear and $\mu_{1}$-measurable.

Proof :
Let us suppose first that $T$ is unitary and surjective.
Let $j_{E_{2}}$ a continuous linear injection of $E_{2}$ into $\mathbb{R}^{\mathbb{N}}$ mapping $\mu_{2}$ onto $\gamma_{1}^{\otimes \mathbb{N}}$. To this injection we associate an orthonormal basis of $H_{2}$, namely $f_{n}: f_{n}=j_{E_{2}}^{-1}\left(e_{n}\right)$, where $\left(e_{n}\right)_{n}$ denotes the canonical basis of $\ell^{2}$.
Let $g_{n}=T^{-1}\left(f_{n}\right)$ : by the hypothesis, $\left(g_{n}\right)$ is an orthonormal basis of $H_{1}$ which defines a (equivalence class of) linear $\mu_{1}$-measurable injection of $E_{1}$ into $\mathbb{R}^{\mathbb{N}}$, let $j_{E_{1}}$, and $j_{E_{1}}$ maps $\mu_{1}$ onto $\gamma_{1}^{\otimes \mathbb{N}}$ but it is neither continuous, nor injective in general.

Let $G=j_{E_{1}}\left(E_{1}\right) \cap j_{E_{2}}\left(E_{2}\right) \subset \mathbb{R}^{\mathbb{N}} . G$ is a vector subspace of $\gamma_{1}^{\otimes \mathbb{N}}$-probability equal to one. Now let us consider

$$
j_{E_{2}}^{-1} \circ j_{E_{1}}: \quad j_{E_{1}}^{-1}(G) \rightarrow E_{2}
$$

It is a linear map, extending $T$.
Now we can extend

$$
j_{E_{2}}^{-1} \circ j_{E_{1}} \text { to } E_{1}
$$

and obtain a linear $\mu_{1}$-measurable extension of $T$ to $E_{1}$. It is clear that :

$$
\widetilde{T}\left(\mu_{1}\right)=\mu_{2}
$$

## In the general case :

We may clearly suppose that $\|T\| \leq 1$.
Let $E_{1} \times E_{2}$ with the Gaussian probability $\mu_{1} \otimes \mu_{2}$, whose Cameron-Martin space is equal to $H_{1} \otimes H_{2}$.
Let

$$
\begin{array}{ll}
A=\sqrt{I_{H_{2}}-T T^{*}}, & A: H_{2} \rightarrow H_{2} \\
B=\sqrt{I_{H_{1}}-T^{*} T}, & B: H_{1} \rightarrow H_{1}
\end{array}
$$

I claim that :

$$
T B=A T \quad \text { and } \quad T^{*} A=B T^{*}
$$

To see this fact, it is sufficient to consider the binomial expansion of

$$
\left(I_{H_{2}}-T T^{*}\right)^{\frac{1}{2}} \quad \text { and } \quad\left(I_{H_{1}}-T^{*} T\right)^{\frac{1}{2}}
$$

Moreover : let $U: H_{1} \otimes H_{2} \rightarrow H_{1} \otimes H_{2}$ defined by :

$$
U(x, y):=\left(-B x+T^{*} y, T x+A y\right)
$$

taking account of the commutation relations above, we can easily verify that $U$ is unitary surjective, with the inverse equal to $U$.
By the first part of the proof there exists an extension

$$
\widetilde{U}: E_{1} \times E_{2} \rightarrow E_{1} \times E_{2}
$$

linear and $\mu_{1} \otimes \mu_{2}$-measurable.

Next define :

$$
v(x, y)=\operatorname{proj}_{E_{2}}\left[\frac{\tilde{U}(x, y)+\widetilde{U}(x,-y)}{2}\right]
$$

$v$ is linear and $\mu_{1} \otimes \mu_{2}$-measurable from $E_{1} \times E_{2}$ into $E_{2}$ extending T. Moreover :

$$
v(x, y)=v(x, 0), \quad \forall x
$$

By Lemma 3, $v(\cdot, 0)$ is $\mu_{1}$-measurable and : $\widetilde{T}(\cdot)=v(\cdot, 0)$ satisfies the property of the theorem.

- Q.E.D.-
$\widetilde{T}$ is called a " $\mu_{1}$-essential extension of $T$ ".
Now we give a particular case which will be useful in the sequel.
THEOREM 7 : Let us suppose that, in addition with the hypothesis of Theorem $6, T$ is Hilbert-Schmidt. Then $\tilde{T}$ has $\mu_{1}$-almost surely its values in $H_{2}$ and we have :

$$
\int_{E_{1}}\|\tilde{T} x\|_{H_{2}}^{2} \mu_{1}(d x)<\infty
$$

Conversely, if $\widetilde{T}$ is a linear $\mu_{1}$-measurable map from $E_{1}$ into $H_{2}$ satisfying the above condition, its restriction $T$ to $H_{1}$ is Hilbert-Schmidt. We have the equality :

$$
\|T\|_{H, S}^{2}=\int_{E_{1}}\|\tilde{T} x\|_{H_{2}}^{2} \mu_{1}(d x)
$$

## Proof :

Let $T u=\sum \lambda_{n}\left(u, e_{n}\right)_{H_{1}} f_{n}$ be a spectral decomposition of $T, \sum \lambda_{n}^{2}<\infty,\left(e_{n}\right)$ and $\left(f_{n}\right)$ be orthonormal basis in $H_{1}$ and $H_{2}$ respectively. Let us define :

$$
\widetilde{T}(x)=\sum \lambda_{n} \widetilde{e}_{n}(x) f_{n}
$$

$\left(\widetilde{e}_{n} \in\left(E_{1}^{\prime}\right)_{2}\right.$ is such that $\left.S\left(\widetilde{e}_{n}\right)=e_{n}\right)$.
This series converges in $L^{2}\left(E_{1}, \mu_{1}, H_{2}\right)$ and $\widetilde{T}$ is an extension of $T$.

Moreover :

$$
\begin{aligned}
\int_{E_{1}}\|\tilde{T} x\|_{H_{2}}^{2} \mu_{1}(d x) & =\sum \lambda_{n}^{2} \\
& =\|T\|_{H, S}^{2} .
\end{aligned}
$$

Conversely let $\widetilde{T}: E_{1} \rightarrow H_{2}$ linear, $\mu_{1}$-measurable with

$$
\int_{E_{1}}\|\tilde{T} x\|_{H_{2}}^{2} \mu_{1}(d x)<\infty
$$

$\mu_{2}=\widetilde{T}\left(\mu_{1}\right)$ is a Gaussian measure on $H_{2}$ with reproducing kernel Hilbert space $T\left(H_{1}\right)$ where $T=\widetilde{T}_{\mid H_{1}}$.
By Minlos' theorem :

$$
T: H_{1} \rightarrow H_{2} \quad \text { is Hilbert-Schmidt. }
$$

- Q.E.D.-

Remark : every $\mu_{1}$-measurable linear function from $E_{1}$ into the Hilbert $H_{2}$ is an almost sure limit of linear continuous mappings from $E_{1}$ into $H_{2}$.

Actually if $H_{1}$ denotes the Cameron-Martin space of $\mu_{1}$ and if $i: H_{1} \rightarrow E_{1}$ is the canonical injection, then $\tilde{T} \circ i: H_{1} \rightarrow H_{2}$ is Hilbert-Schmidt :

$$
\widetilde{T} \circ i(u)=\lim _{n \rightarrow \infty} \sum_{k=0}^{n} \lambda_{k}\left(e_{k}, u\right) f_{k}
$$

Therefore $\widetilde{T}(\bullet)$ is a $L^{2}$-limit of mappings :

$$
\sum_{k=0}^{n} \lambda_{k} \tilde{e}_{k}(\cdot) f_{k}
$$

But such a mapping is a $L^{2}$-limit of continuous mappings with finite rank, since the continuous linear functionals are dense in $\left(E_{1}\right)_{2}^{\prime}$.

Therefore each linear measurable function from $E_{1}$ into $H_{2}$ is a $L^{2}$-limit of continuous linear mappings.

By extracting some subsequence : we have the announced result.

## 4 - Some applications

## a). Quantization of a linear measurable operator

Let $(E, H, \mu)$ be a Wiener space, let $T: H \rightarrow H$ linear, continuous with $\|T\| \leq 1$ and let $B=\left(I-T^{*} T\right)^{\frac{1}{2}}$. The $\mu$-essential extensions of $T^{*}$ and $B$ will be denoted by $\widetilde{T}^{*}$ and by $\widetilde{B}$ respectively ( $T^{*}$ is the adjoint of $T$ ). Let us set

$$
\Gamma(T) f(x)=\int_{E} f\left(T^{*} x+\widetilde{B} y\right) \mu(d y)
$$

where $f$ is a Borelian (or $\mu$-measurable) positive function.
Then $\Gamma(T) f$ does not depend from the extensions chosen. Moreover :

$$
\forall p \geq 1, \quad \Gamma(T): L^{p}(E, \mu) \rightarrow L^{p}(E, \mu) \quad \text { is contracting } .
$$

Actually if $f$ is Borelian, $\Gamma(T) f$ is well defined since

$$
(x, y) \rightsquigarrow \widetilde{T}^{*} x+\widetilde{B} y \quad \text { is } \mu \otimes \mu \text {-measurable }
$$

and we have seen that :

$$
\text { for every } \mathrm{x}: \quad y \rightsquigarrow \widetilde{T}^{*} x+\widetilde{B} y \text { is } \mu \text {-measurable. }
$$

Moreover the functions

$$
(x, y) \rightsquigarrow f\left(\widetilde{T}^{*} x+\widetilde{B} y\right)
$$

and

$$
x \leadsto f(x)
$$

have the same law and therefore :

$$
N_{p}(\Gamma(T) f) \leq N_{p}(f), \quad 1 \leq p \leq \infty
$$

b) . Linear measurable functions on $E=\mathcal{C}_{0}([0,1])$ with the Wiener-measure, denoted $P$

The dual of $E$ is isomorphic to the space $\mathcal{M}([0,1])$ of real Borel measures on $] 0,1]$. If $\left.\left.\varphi \in L^{2}(] 0,1\right], d t\right)$, the random variable :

$$
\int_{0}^{1} \varphi(t) d \omega(t) \quad \text { (Wiener-Itô integral) }
$$

is a $P$-measurable linear form. Actually it is a $L^{2}$-limit of random variables from the form

$$
\sum \varphi\left(t_{i}\right)\left[\omega_{t_{i+1}}-\omega_{t_{i}}\right]
$$

which are linear and continuous.

We shall see that all $P$-measurable functionals $L$ are of this form.

## First let $L$ be a continuous linear form :

Then there exists an unique (signed) measure $\mu$ on $] 0,1]$ such that :

$$
\begin{aligned}
\left.\left.\forall \omega \in \mathcal{C}_{0}(] 0,1\right]\right), \quad L(\omega) & =\int_{[0,1]} \omega(t) \mu(d t) \\
& =\int_{[0,1]} \omega(t) d F_{\mu}(t)
\end{aligned}
$$

But by the stochastic integration by parts formula, we obtain :

$$
\int_{[0,1]} \omega(T) d F_{\mu}(t)=-\int_{0}^{1} F_{\mu}(t) d \omega(t)+\left[F_{\mu}(t) \omega(t)\right]_{0}^{1}
$$

(The integral in the right member is a Wiener-Itô integral).
Therefore

$$
\left.\left.L(\omega)=\int_{0}^{1} \mu(] t, 1\right]\right) d \omega(t)
$$

## In the general case :

We approach $L$ (in the $E_{2}^{\prime}(P)$-topology) by linear continuous forms and it suffices to notice that the functions

$$
s \leadsto \mu(] s, 1]), \quad(\mu: \text { real measure on }] 0,1])
$$

are dense in $\left.L^{2}([0,1], d t]\right)$.

Remark : In the case of the Brownian motion indexed by $\mathbb{R}_{+}$ we can prove by absolutely the same way that every linear measurable form on $\mathcal{C}_{0}([0, \infty[, \mathbb{R})$, equipped with the Wiener measure is the Itô integral of an (unique) element of $L^{2}([0, \infty[, d t)$. The converse is also true .
c) . Infinite dimensional Brownian motion

Let $\left(\mathbb{B}_{t}\right)_{t \in \mathbb{R}^{+}}=\left(\left(\mathbb{B}_{n}(t)\right)_{n \in \mathbb{N}}\right)_{t \in \mathbb{R}^{+}}$be an infinite-dimensional Brownian motion with parameter set $\mathbb{R}_{+}$or $[0,1]$.

Then it corresponds to it the following (abstract) Wiener space :

$$
E=\mathcal{C}_{0}\left(\left[0, \infty\left[, \mathbb{R}^{\mathbb{N}}\right), \quad H=W^{1,2}\left(\left[0, \infty\left[, l^{2}\right), \quad P=\left(\mu_{1}\right)^{\otimes \mathbb{N}}\right.\right.\right.\right.
$$

where $\mu_{1}$ is the ordinary Wiener measure.
As above we see that every linear measurable functional is a stochastic integral with respect to $\left(\mathbb{B}_{t}\right)$ of an element of $L^{2}\left(\left[0, \infty\left[, l^{2}\right)\right.\right.$. This means :

$$
L(\cdot)=\int<f(t), d \mathbb{B}_{t}>:=\sum_{n} \int f_{n}(t) d \mathbb{B}_{n}(t)
$$

The conclusion remains true for a Brownian motion with values in the abstract Wiener space ( $F, H, \mu$ ) (see chapter one) : every linear measurable functional is the stochastic integral of an element of $L^{2}([0, \infty[, H)$.

## d) . Expansions in "Fourier series"

Let $(E, H, \mu)$ be an abstract Wiener space and $j: E \rightarrow \mathbb{R}^{\mathbb{N}}$ a continuous linear injection from $E$ into $\mathbb{R}^{\mathbb{N}}$ such that $j(\mu)=\gamma_{1}^{\otimes \mathbb{N}}$.
Let $\widetilde{e}_{n}$ the (continuous) linear form on $E$ defined as

$$
\widetilde{e}_{n}(x):=(j(x))_{n}
$$

and let $e_{n}=S\left(\widetilde{e}_{n}\right)$, the barycenter ; we have :

$$
e_{n} \in H
$$

I claim that :

$$
x=\sum_{n=0}^{\infty} \tilde{e}_{n}(x) \cdot e_{n}, \quad \text { almost surely }
$$

Actually let $\left(f_{n}\right)$ the canonical "basis" of $\mathbb{R}^{\mathbb{N}}$.
Then

$$
f_{n}=j\left(e_{n}\right)
$$

In $\mathbb{R}^{\mathbb{N}}$ we have the expansion :

$$
y=\sum y_{n} f_{n}
$$

(if $\left.y=\left(y_{n}\right)_{n}\right)$

Now let $K_{\varepsilon}$ be a compact set in $\mathbb{R}^{\mathbb{N}}$, contained in $j(E)$ and such that :

$$
\gamma_{1}^{\otimes \mathbb{N}}\left(K_{\varepsilon}\right) \geq 1-\varepsilon
$$

and let

$$
L_{\varepsilon}=j^{-1}\left(K_{\varepsilon}\right)
$$

Then $j$ defines an homeomorphim from $L_{\varepsilon}$ onto $K_{\varepsilon}$, and we have the equivalence :

$$
\left(x \in K_{\varepsilon}, \quad \sum \widetilde{e}_{n}(x) e_{n}=x\right) \Longleftrightarrow\left(y \in L_{\varepsilon}, \quad y=\sum y_{n} f_{n}\right)
$$

Therefore :

$$
x=\sum \widetilde{e}_{n}(x) e_{n}, \quad \text { for every } x \text { in } K_{\varepsilon}
$$

Finally on $\bigcup_{m \geq 1} K_{\frac{1}{m}}$, ( a Borel set carrying $\mu$ ), we have :

$$
x=\sum \widetilde{e}_{n}(x) e_{n}
$$

and the result is proven for $\left(e_{n}\right)=\left(j^{-1}\left(f_{n}\right)\right)$.
Now let $\left(\widetilde{e}_{n}^{1}\right)$ be an orthonormal basis of $E_{2}^{\prime}(\mu)$, not necessarily contained in $E^{\prime}$. We set as before

$$
e_{n}^{1}=S\left(\tilde{e}_{n}^{1}\right)
$$

Let $\widetilde{e}_{n}$ and $\left(e_{n}\right)$ as above, $\left(e_{n}=j^{-1}\left(f_{n}\right)\right)$. To this change of basis, we associate an unitary operator $A: H \rightarrow H$. Let $\tilde{A}$ a linear $\mu$-measurable extension of $A, \tilde{A}: E \rightarrow E$; we have $\widetilde{A}(\mu)=\mu$.
Let $K_{\varepsilon} \subset E$ be a compact such that :

- $\mu\left(K_{\varepsilon}\right) \geq 1-\varepsilon$
- $\sum \widetilde{e}_{n}(x) e_{n}=x, \quad \forall x \in K_{\varepsilon}$
- the restriction of $\widetilde{A}$ to $K_{\varepsilon}$ is continuous.

Then :

$$
x \in K_{\varepsilon} \Rightarrow \sum \widetilde{e}_{n}(x) \widetilde{A}\left(e_{n}\right)=\tilde{A}(x)
$$

Therefore

$$
\sum \widetilde{e}_{n}(x) e_{n}^{1}=\widetilde{A}(x) \quad \text { for every } x \in K_{\varepsilon}
$$

But

$$
\tilde{e}_{n}(x)=\tilde{e}_{n}^{1}(\tilde{A} x) \quad \text { almost surely }
$$

In fact the two measurable forms on $E$ :

$$
x \leadsto \widetilde{e}_{n}(x) \quad \text { and } \quad x \leadsto\left(A e_{n}, \widetilde{A} x\right)=\widetilde{A e_{n}}(\tilde{A} x)
$$

coincide on $H$ since $A$ is unitary, therefore coincide almost everywhere on $E$. We can therefore suppose that :

$$
\tilde{e}_{n}(x)=\tilde{e}_{n}^{1}(\tilde{A} x), \quad \forall x, \forall n
$$

and this means that:

$$
\sum \tilde{e}_{n}^{1}(y) e_{n}^{1} \longrightarrow y \quad \text { for } y \in \tilde{A}\left(K_{\varepsilon}\right)
$$

Since $\widetilde{A}\left(K_{\varepsilon}\right)$ carries $\mu$ up to $\varepsilon$, we have the announced result.

