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# INFINITELY DIVISIBLE MEASURES ON <br> MULTI-HILBERTIEN SPACES AND 

A LÉVY-ITO DECOMPOSITION

By

Milan J. Merkle

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# ABSTRACT <br> INFINITELY DIVISIBLE MEASURES ON MULTI-HILBERTIEN SPACES AND A LEVY-ITO DECOMPOSITION 

By

Milan J. Merkle

In this work, we give a representation of infinitely divisible (ID) laws on duals of multi-Hilbertien spaces and discuss the convergence. These results give a unified approach to the existant work on infinitely dimensional Hilbert spaces and on nuclear spaces. The convergence of ID laws can be used to prove the weak convergence of homogeneous processes with independent increments. This is applied to a problem from Neurobiology, and the results obtained are generalization and improvement of a recent work of G. Kallianpur [16]. The last chapter is devoted to the processes with independent increments on duals of multi-Hilbertien spaces. Lévy-Ito decomposition on Hilbert spaces is obtained.

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Special thanks are due to Dr. V. Mandrekar, whose ideas and suggestions were of greatest value for this work.

Many memorable hours are spent with my friends, graduate assistants; among other things, drinking American coffee and eating Indian food is what makes life exciting.

## TABLE OF CONTENTS

0. INTRODUCTION ..... 1
1. PRELIMINARIES ..... 3
1.1. Vector spaces and seminorms ..... 3
1.2. Operators on Hilbert spaces ..... 3
1.3. Multi-Hilbertien spaces ..... 5
1.4. Dual spaces ..... 9
1.5. Measures on dual spaces ..... 9
2. REPRESENTATION OF ID LAWS ..... 15
2.1. I $(\tau)$-topology and Bochner's theorem ..... 15
2.2. Continuity theorem ..... 18
2.3. Infinitely divisible probability measures ..... 23
3. CONVERGENCE OF ID LAWS AND HOMOGENEOUS PROCESSES WITH INDEPENDENT INCREMENTS ..... 33
3.1. Convergence of ID laws ..... 33
3.2. Homogeneous processes with independent increments ..... 34
3.3. An example ..... 36
4. PROCESSES WITH INDEPENDENT INCREMENTS ..... 41
4.1. Properties of paths ..... 41
4.2. A decomposition of E'-valued processes ..... 44
4.3. A Lévy-Ito decomposition ..... 48
REFERENCES ..... 59

## 0. INTRODUCTION

The infinitely divisible laws (ID) in the context of infinitedimensional Hilbert spaces were studied by S.R.S. Varadhan ([23]). For the case of nuclear spaces, a Lévy-Khinchine representation was given by Fernique [6]. Our purpose here is to give a representation of ID laws on a multi-Hilbertien space and discuss the convergence. These results give unified approach to Varadhan's and Fernique's work and extend Fernique's work to include convergence of ID laws. Such a unified approach for Bochner theorem for multi-Hilbertien spaces is recently given by K. Ito [14]. Our results are based on arguments of Fernique for the representation and Levy continuity theorem for the multi-Hilbertien spaces.

We note that the method of Fernique for the Lévy continuity theorem in nuclear spaces fails in the multi-Hilbertien spaces due to an example of Sazonov (cf. 2.1.2).

The convergence of ID laws can be used to prove the weak convergence of additive (homogeneous independent increments) processes. This is shown in the context of an example from Neurobiology. This work includes and improves recent work of G. Kallianpur [16]. In fact, our work is applicable to wider class of examples.

The last chapter is devoted to the Lévy-Ito decomposition of infinitedimensional processes with independent increments (PII). We are unable at this point to prove this theorem for general multi-Hilbertien spaces. In fact, as the reader can see, the results seem to need some work even for the Hilbert space case. However, for the additive processes, such a decomposition follows from the relation of convolution semigroups to ID laws. This, along with the convergence of the related processes, is studied in Chapter 3. Clearly, the results there allow us to extend the
work of Kallianpur to multi-neuronal models.
The convergence of PII in this context remains an open problem. But our decomposition is the basic technique needed for this study. In future work we intend to study this problem as a generalization of a recent work of Jacod (ZWV 63, 109-136, 1983).

## 1. PRELIMINARIES

### 1.1. Vector Spaces and Seminorms

$R(N)$ will denote the set of real (natural) numbers. Throughout this work the phrase "vector space" will mean a real vector space.

Let $E$ be a vector space. A real-valued function $p$ defined on $E$ is called a seminorm if for all $x, x_{1}, x_{2} \in E$ and $a \in R$ :

$$
\begin{array}{ll}
\text { (i) } & p(x) \geq 0 \\
\text { (ii) } & p(a x)=|a| p(x) \\
\text { (iii) } & p\left(x_{1}+x_{2}\right) \leq p\left(x_{1}\right)+p\left(x_{2}\right), \text { and }  \tag{iii}\\
\text { (iv) } & p(x)>0 \text { for some } x .
\end{array}
$$

A seminorm $p$ on $E$ is called a Hilbertien seminorm if

$$
\left(\forall x_{1}, x_{2} \in E\right) \quad p^{2}\left(x_{1}+x_{2}\right)+p^{2}\left(x_{1}-x_{2}\right)=2\left(p^{2}\left(x_{1}\right)+p^{2}\left(x_{2}\right)\right)
$$

For a Hilbertien seminorm $p$ we define a symmetric bilinear form $p(\cdot, \cdot)$ on $E$ by

$$
p\left(x_{1}, x_{2}\right)=\frac{1}{4}\left(p^{2}\left(x_{1}+x_{2}\right)-p^{2}\left(x_{1}-x_{2}\right)\right)
$$

$p\left(x_{1}, x_{2}\right)$ is called the inner product corresponding to $p$.
A seminorm $p$ is said to be separable if there is a countable set $D \subset E$ such that for each $x \in E$ and $\varepsilon$ positive, there exists a $d$ in $D$ such that $p(x-d) \leq \varepsilon$.

### 1.2 Operators on Hilbert Spaces

Let $H_{i}(i=1,2)$ be separable real Hilbert spaces with norms $\|\cdot\|_{i}$
and inner products $\langle\cdot, \cdot\rangle_{i}$. Let $A$ be a linear mapping from $H_{1}$ to $\mathrm{H}_{2}$.
$A$ is a compact operator if for every $x$ in $H_{1}$ :
(1)

$$
A x=\sum_{n} t_{n}<x, e_{n}>h_{n}
$$

where $0<t_{n} \downarrow 0$, and $e_{n}, h_{n}$ are orthonormal sets in $H_{1}$ and $H_{2}$ respectively.
$A$ is a Hilbert-Schmidt operator if (1) holds and
(2)

$$
\sum_{n} t^{2}<\infty
$$

$A$ is a nuclear operator if (1) holds and

$$
\begin{equation*}
\sum_{n} t_{n}<\infty \tag{3}
\end{equation*}
$$

Let $H_{1}=H_{2}=H$. Operator $A$ is said to be a trace class operator if

$$
\begin{equation*}
\sum_{n}<A e_{n}, e_{n}>1<\infty \tag{4}
\end{equation*}
$$

for any orthonormal basis $\left\{e_{n}\right\}$ in $H$.
For a linear mapping $A: H \rightarrow H$ we define the adjoint mapping $A^{*}$ by

$$
\begin{equation*}
\left\langle A^{*} x_{1}, x_{2}\right\rangle=\left\langle x_{1}, A^{*} x_{2}\right\rangle \tag{5}
\end{equation*}
$$

We say that $A$ is self-adjoint if $A=A^{*}$.

A is said to be positive (non-negative definite) if

$$
\begin{equation*}
\sum_{i, j} a_{i} \bar{a}_{j}<x_{i}, A x_{j}>\geq 0 \tag{6}
\end{equation*}
$$

for every finite set of complex numbers $\left\{a_{i}\right\}_{i=1}^{n}$ and $\left\{x_{i}\right\}_{i=1}^{n}$ a finite subset of $H$.

In the following theorem we give a summary of some well-known facts (see, for example [9], Chapter 1.)

Theorem 1. (i) If $A$ and $B$ are Hilbert-Schmidt operators, then $A B$
is a nuclear operator.
Conversely, every nuclear operator is a product of two Hilbert-Schmidt operators.
(ii) A compact positive operator $A: H \rightarrow H$ is nuclear if and only if it is a trace operator.
(iii) If $A$ is nuclear (or Hilbert-Schmidt) operator, then $A^{*}$ is an operator of the same type.
(iv) If $A: H \rightarrow H$ is a positive nuclear operator then $A^{\frac{1}{2}}$ is positive Hilbert-Schmidt operator.
(v) A positive compact operator $A: H \rightarrow H$ has the representation.

$$
A x=\sum t_{n}<x, e_{n}>e_{n}
$$

where $e_{n}$ is an orthonormal basis for $H$ consisting of eigenvectors of $A$ and $t_{n}, 0<t_{n}+0$ are eigenvalues of $A, A e_{n}=t_{n} e_{n}$.
(vi) Let $H_{1}$ and $H_{2}$ be Hilbert spaces and $A: H_{1} \rightarrow H_{2}$ be a linear mapping. $A$ is a Hilbert-Schmidt operator if and only if

$$
\sum_{n}\left\|A e_{n}\right\|_{2}^{2}<\infty
$$

for at least one (equivalently: for all) orthonormal basis $\left\{e_{n}\right\}$ in $H_{1}$. $A$ is nuclear operator if and only if

$$
\sum_{n}\left|<A e_{n}, e_{n}>2\right|<\infty .
$$

In the following section we describe multi-Hilbertien spaces, following Ito [14].

### 1.3. Multi-Hilbertien Spaces

Let $E$ be a vector space.
The family of all separable Hilbertien seminorms on $E$ will be denoted
by HSN.

It is easy to see that if $p, p_{1}, \ldots, p_{n} \in H S N$ then for every positive $c, c p \in H S N$ and $\left(p_{1}^{2}+\ldots+p_{n}^{2}\right)^{\frac{1}{2}} \in H S N$.

For $p, q \in H S N$ we define the relation $\ll p \geqslant q$ if and only if $p(x) \leq c q(x)$ for some $c \in R$ and all $x \in E$.

Let $E_{p}$ denote $E$ with topology given by the seminorm $p$.
Let $\operatorname{ker} p=\{x \in E: p(x)=0\}$. Then $E_{p} /$ ker $p$ is a pre-Hilbert space; its completion $\bar{E}_{p}$ is a Hilbert space.

If $p \geqslant q$ then the identity map $E \rightarrow E$ can be extended to a continuous linear operator $i_{p, q}: \bar{E}_{2} \rightarrow \bar{E}_{p}$. Notice that $i_{p, q}(x)=x$ on $E$, and the indices $p, q$ are just pointing the topology for completion.

Define the relation $<{ }_{H S} \quad: \quad p<q$ if and only if $p<q$ and $i_{p, q}$ is a Hilbert-Schmidt operator.

In view of Theorem 3.1.(vi), $p<{ }_{H S} q$ if and only if

$$
\sum p^{2}\left(e_{n}\right)<\infty
$$

where $\left\{e_{n}\right\}$ is an orthonormal basis in $\bar{E}_{q}$.
Let $P=\left\{p_{d}\right\}_{j \in J}$ be a subfamily of HSN on $E$ such that
(i) If $p \in P$ and $p<q$ or $q<p$ then $q \in P$
(ii) If $p, q \in P, p\left(x_{n}\right) \rightarrow 0, x_{n}$ is $q$-Cauchy then $q\left(x_{n}\right) \rightarrow 0$.

Define a locally convex topology $\tau$ on $E$ by its neighborhood basis at zero:
(1)

$$
\begin{array}{r}
\left\{x \in E: p_{j_{1}}(x)<1, \ldots, p_{j_{k}}(x)<1, j_{1}, \ldots, j_{k} \in J\right\} \\
k=1,2, \ldots .
\end{array}
$$

The topology $\tau$ given by (1) is called a multi-Hilbertien topology generated by $P$.

Space E with a multi-Hilbertien topology $\tau,(E, \tau)$ is called a multiHilbertien space.

If a multi-Hilbertien topology $\tau$ is generated by $P=\left\{p_{j}\right\}_{j \in J}$ we will often write $\tau=\left\{p_{j}\right\}_{j \in J}$.

If $\tau=\left\{p_{j}\right\}_{j \in J}$ where $J$ is a countable set. we say that $(E, \tau)$ is a countably Hilbert space. In this case, without loss of generality, we may assume

$$
\begin{equation*}
p_{1}<p_{2}<p_{3}<\cdots \tag{2}
\end{equation*}
$$

For, if not, then the topology given by $a_{j}=\left(p_{1}^{2}+p_{2}^{2}+\ldots+p_{j}^{2}\right)^{\frac{1}{2}}$, $j=1,2, \ldots$ coincide with the topology given by $\left\{p_{j}\right\}$ and $\left\{q_{j}\right\}$ satisfies condition (2).

In a countably Hilbert space we will use notations $E_{n}, \bar{E}_{n}, E_{n}^{\prime}$, rather than $E_{P_{n}}$, etc. ...

If the topology $\tau$ is given by countably many norms $\left\{p_{n}\right\}$ (i.e. ker $p_{n}=$ $\{0\}$ ) , then we can define a metric $d$ by

$$
\begin{equation*}
d\left(x_{1}, x_{2}\right)=\sum_{n} 2^{-n} \frac{p_{n}\left(x_{1}-x_{2}\right)}{1+p_{n}\left(x_{1}-x_{2}\right)} \tag{3}
\end{equation*}
$$

$E$ is said to be complete if ( $E, d$ ) is complete.
$E$ is complete if and only if ([27])

$$
\begin{equation*}
E=\hat{n}_{n} \bar{E}_{p_{n}} \tag{4}
\end{equation*}
$$

We say that ( $E, \tau$ ) is a nuclear space if for every $p \in \tau$ there is a $a \in \tau$ such
that $p<S_{S}$. Notice that, by Theorem 3.1.(i), $E$ is nuclear if and only if for every $p \in \tau$ there is a $q \in \tau$ such that $i_{p, q}$ is nuclear operator. A nuclear space in which condition (4) holds is called a nuclear Fréchet space.

Suppose that $\left\{e_{j}^{n}\right\}_{j=1}^{\infty}$ is a complete orthonormal basis is $\bar{E}_{n}$, $n=1,2, \ldots$.

Then $E$ is a nuclear space if and only if for every $n=1,2, \ldots$ there is an $m>n$ such that (Theorem 3.1.(vi)).

$$
\begin{equation*}
\sum_{j} p_{m}^{2}\left(e_{j}^{n}\right)<\infty \tag{5}
\end{equation*}
$$

Example. Let $C^{\infty}(R)$ be the set of all infinitely differentiable functions defined on $R$. The Schwartz space $s$ is a locally convex subspace of $C^{\infty}(R)$, topology defined by norms $q_{n}(f)=\max _{0 \leq k \leq n} \sup _{t}\left|\left(1+t^{2}\right)^{n_{f}}{ }^{(n)}(t)\right|$. It can be shown that this topology is equivalent to the one given by Hilbertien norms:

$$
p_{n}(f)=\left(\sum_{k=0}^{n} \int_{R}\left(1+t^{2}\right)^{n}\left|f^{(k)}(t)\right|^{2} d t\right)^{\frac{1}{2}}
$$

Let $\tau=\left\{p_{n}\right\}_{n=1}^{\infty}$. Then $(s, \tau)$ is a nuclear space, with $E_{n}=\left\{f \in C^{\infty}(R): p_{n}(f)<\infty\right\}$ and the inner products:

$$
p_{n}\left(f_{1}, f_{2}\right)=\sum_{k=0}^{n} \int_{R}\left(1+t^{2}\right)^{n} f_{1}^{(k)}(t) f_{2}^{(k)}(t) d t
$$

Remark. A separable Hilbert space $H$ is nuclear if and only if $\operatorname{dim} H<\infty$.

Let now $(E, \tau)$ be a multihilbertian space, $\tau=\left\{p_{j}\right\}_{j \in J}$. We say that a set $B \subset E$ is bounded if for each $j \in J$ the set of numbers

$$
\left\{p_{j}(x): x \in B\right\}
$$

is bounded.

### 1.4 Dual Spaces

Let ( $E, \tau$ ) be a multi-Hilbertien space. The set of all ( $\tau-$ ) continuous linear functionals on $E$ is denoted by $E_{\tau}^{\prime}$ (or $E^{\prime}$ if there is no confusion about the topology on $E$ ).

If $F \in E^{\prime}, x \in E$, then $F$ evaluated at $x$ is denoted by $\langle F, x\rangle$. In this dual notation it is understood that for fixed $F$ (or $x$ ), $<\mathrm{F}, \mathrm{x}>$ is a function of $x$ (or F). Also, $\langle F, x\rangle$ separates points in both $E$ and $E^{\prime}$. The strong topology on $E^{\prime}$ is given by seminorms

$$
\|F\|_{B}=\sup _{x \in B}|\langle F, x\rangle|, B \text {-bounded set in } E .
$$

Let $E_{p}^{\prime}$ denote the topological dual of $E_{p} . E_{p}^{\prime}$ is a Hilbert space with the norm $\bar{p}(F)=\sup _{p(x) \leq 1}|<F, x>|$.

We have $E^{\prime}=\bigcup_{p} E_{p}^{\prime}$.
In the strong topology, $E^{\prime}$ is an inductive limit of $\left\{E_{p}^{\prime}\right\}_{p \in \tau}$ [14].
If $\tau$ is countable, the set $A$ is bounded (compact) in $E^{\prime}$ if and only if it is bounded (compact) in some $E_{p}^{\prime}, p \in \tau$ [8].
1.5 Measures on dual spaces

Let $E^{\prime}$ be a topological dual of a topological vector space $E$.

Let $A$ be a given Borel set in $R^{n}$. Let
(1)

$$
Z=\left\{F \in E^{\prime}:\left(\left\langle F, x_{1}\right\rangle,\left\langle F, x_{2}\right\rangle, \ldots,\left\langle F, x_{n}\right\rangle\right) \in A\right\}
$$

Set $Z$ defined by (1) for some $x_{1}, \ldots, x_{n} \in E$ is called the cylinder set with base $A$ and generating elements $x_{1}, \ldots, x_{n}$.

Another approach to defining a cylinder set is the following:

Let $Y$ be a finitely dimensional subspace of $E$. Let $Y^{0}$ denote the annihilator of $Y$, i.e., the set of all $F \in E^{\prime}$ such that

$$
\langle F, x\rangle=0 \quad \text { for } x \in Y
$$

Consider the factor space $E^{\prime} / Y^{0}$. It is isomorphic to $Y^{\prime}$, thus finitely dimensional. Let $A \subset E^{\prime} / Y^{0}$ be a Borel set. The set of all $F \in E^{\prime}$ which are carried into elements of $A$ by the natural mapping $E^{\prime} \rightarrow E^{\prime} / Y^{0}$ is called the cylindar set with base $A$ and generating subspace $Y^{\circ}$.

These two definition are equivalent and define the same object ([9]). Note that a cylinder set may have more than one representation in terms of base and generating subspace.

It is easy to see that the cylinder sets form an algebra of sets-cylinder algebra.

By a cylinder set measure in $E^{\prime}$ we mean a nonnegative function $M$ defined on the cylinder algebra with the following properties:
(i) If $Z=\bigcup_{i=1}^{\infty} Z_{i}$ where $Z_{i} \cap Z_{j}=\emptyset$ if $i \neq j$, and all $Z_{i}$ are generated by the same set $x_{1}, \ldots, x_{n}$, then $M(Z)=\sum_{i=1}^{\infty} M\left(Z_{i}\right)$. (ii) For any cylinder set $Z, M(Z)=\inf M(U)$, where $U$ runs through all open cylinder sets containing $Z$.

Let $M$ be a cylinder set measure. For $Y \subset E$ a finite dimensional subspace we define

$$
\begin{equation*}
M_{\gamma}(A)=M(Z) \tag{2}
\end{equation*}
$$

where $A$ is a Borel set in $E^{\prime} / Y^{\circ}$ and $Z$ is the cylinder with base $A$
and generating subspace $Y^{0}$. So, (2) defines a Borel measure on $E^{\prime} / Y^{\circ}$ which is regular, i.e.

$$
\begin{equation*}
M_{Y}(A)=\inf _{U} M(U) \tag{3}
\end{equation*}
$$

where $U$ runs through all open sets containing $A$.
Measures $M_{y}$ are compatible in the following sense. Let $\gamma_{1} \subset \gamma_{2}$. Let $T$ be the natural mapping $E^{\prime} / Y_{2}^{0} \rightarrow E^{\prime} / Y_{1}^{0}$. Then

$$
\begin{equation*}
M_{Y_{1}}(A)=M_{Y_{2}}\left(T^{-1}(A)\right) \tag{4}
\end{equation*}
$$

Conversely, if the system of measures $\left\{M_{\gamma}\right\}$ is given, satisfying condition (3) and (4), then there is a unique cylinder set measure $M$ such that (2) holds for every $Y$.

If $M$ is a cylinder set measure, then for disjoint cylinder sets $z_{1}, \ldots, z_{n}$, we have $M\left(\bigcup_{i=1}^{n} z_{i}\right)=\sum_{i=1}^{n} M\left(z_{i}\right)$, which follows from (i) and the fact that for finitely many cylinder set there always exists a common generating set. However, a cylinder set measure may not be countably additive. If it is so, then it can be extended to a measure on cylinder $\sigma$-algebra $C$, which is defined to be the smallest $\sigma$-algebra that contains cylinder algebra.

We say that the measure $M$ is a probability measure if $M\left(E^{\prime}\right)=1$.
Let now $E$ be a countably Hilbert space. Then ([14]) the cylinder $\sigma$-algebra $\&$ coicides with the Borel $\sigma$-algebra generated by the strongly open sets. Every probability measure on $E^{\prime}$ is regular.

Let $(E, \tau)$ be a multi-Hilbertian space.
A probability measure $P$ on ( $E_{\tau}^{\prime}, \mathcal{\ell}$ ) is called separable if there exists a countably Hilbert topology $\tau^{\prime} \subset \tau$ such that

$$
P\left(E^{\prime} \tau^{\prime}\right)=1
$$

A probability measure $P$ on $E^{\prime}$ is called infinitely divisible (ID) if, for every $n \in N, P$ can be represented as $n$-th convolution power of some other probability measure $P_{n}$. In terms of random variables, $X$ is an ID random variable if for every $n \in N$ it can be represented as a sum of $n$ independent identically distributed (iid) random variables. We say that the sequence $M_{n}$ of measures converges weakly to a measure $M \quad\left(M_{n} \Rightarrow M\right) \quad$ if for every continuous bounded real function $f$ defined on $E^{\prime}$ we have, as $n \rightarrow \infty$ :

$$
\int_{E^{\prime}} f(F) d M_{n}(F) \rightarrow \int_{E^{\prime}} f(F) d M(F)
$$

For a measure $M$ on $E^{\prime}$ we define its characteristic functional as a complex-valued function defined on $E$ by

$$
f(x)=\int_{E^{\prime}} e^{i\langle F, x\rangle} d M(F) .
$$

The characteristic functional of an ID random variable will be called an ID characteristic functional

A set $M$ of measures on $E^{\prime}$ is said to be (weakly) relatively compact if every sequence $M_{n}$ in $M$ contains a weakly convergent subsequence. A set $M$ of measures on $E^{\prime}$ is called tight if for every $\varepsilon>0$ there is a compact set $K$ such that $M\left(K^{c}\right)<\varepsilon$ for every $M \in M$. A result in [30] confirms the validity of Prohorov's theorem in $E^{\prime}$, i.e., a sequence $M_{n}$ of separable probability measures is relatively compact if and only if it is tight.

We observe that a sequence of separable measures $M_{n}$ is weakly convergent to a measure $M$ if and only if it is relatively compact and the sequence $f_{n}$ of characteristic functionals $f_{n}$ of $M_{n}$ converges pointwise to the characteristic function $f$ of $M$. This follows from the fact that characteristic functional determines the measure uniquely. Example: (Gaussian measures) Let $C$ be a complex valued bilinear fundtimon defined on $E \times E$, satisfying $C(x, x) \geq 0, C\left(x_{1}, x_{2}\right)=\overline{C\left(x_{2}, x_{1}\right)}$, continuous in both arguments, and non-degenerate $(C(x, x)=0 \Rightarrow x=0)$ Let $Y$ be a $n$-dimensional subspace of $E$. We define a measure $g_{Y}$ :

$$
g_{Y}(A)=\frac{1}{(2 \pi)^{\frac{n}{2}}} \int_{A} \exp \left(-\frac{1}{2} C(y, y)\right) d y
$$

where $d y$ is the Lebesgue measure in $Y$ corresponding to the inner product C.

Finitely dimensional Euclidean space $Y$ with inner product $C$ is isomorphic to $Y^{\prime}$ which is isomorphic to $E / Y^{0}$. Therefore, there is natural isomorphism $T_{Y}$ between $Y$ and $E^{\prime} / Y^{\circ}$. Now define a measure $G_{Y}$ on $E^{\prime} / Y^{\circ}$ by

$$
\begin{equation*}
G_{Y}(B)=g_{Y}\left(A_{Y}^{-1}(B)\right) \tag{2}
\end{equation*}
$$

Now we have a set of finite dimensional measures. It can be shown [9] that (2) defines a compatible set of measures; thus, a cylinder set measure $G$ on $E^{\prime}$ is determined. We call it centered Gaussian cylinder measure.

If $E$ is a nuclear space, every cylinder measure is countably additive, therefore it can be extended to cylinder $\sigma$-algebra.

In the case of general countably Hilbert space, a sufficient condition for countable additivity of $G$ is that for some $n$, the identity map-
ping from $E_{n}$ into $E_{c}$ is Hilbert-Schmidt operator, where $E_{c}$ is $E$ topologized by the norm $C^{\frac{1}{2}}(x, x)$.

The characteristic functional of centered Gaussian measure $G$ defined as above is

$$
f(x)=\exp \left(-\frac{1}{2} C(x, x)\right)
$$

The function $C$ is called covariance and, as we have seen, it uniquely determines a centered Gaussian measure.

If $C$ is not non-degenerate, i.e., if for all $x$ in some linear subspace $X$ we have $C(x, x)=0$, then $C$ is nondegenerate on $E / X$, so a Gaussian measure $G_{1}$ on $X^{0}$ can be constructed following the procedure defined above. Then we define a Gaussian measure on $E^{\prime}$ by

$$
G(A)=G_{1}\left(A \cap X^{O}\right)
$$

Finally, if $G$ is a centered Gaussian measure, $\delta_{F}$ the measure that gives mass 1 to some element $F$ of $E^{\prime}$, then a noncentered Gaussian measure is defined by $\delta_{F} * G$, where * denotes the convolution. The characteristic functional of noncentered Gaussian measure $\delta_{F} * G$ is

$$
f(x)=\exp \left(i<F, x>-\frac{1}{2} C(x, x)\right)
$$

Gaussian measure is an infinitely divisible measure.

## 2. REPRESENTATION OF ID LAWS

### 2.1. I( $\tau)$--Topology and Bochner's Theorem

1. Definition. Let $(E, \tau)$ be a multi-Hilbertien space. We denote by $I(\tau)$ the Hilbert-Schmidt topology induced by all those Hilbertien seminorms which are $<{ }_{H S}$ to some seminorm in $\tau$, i.e.

$$
\begin{equation*}
I(\tau)=\left\{q \in H S N: q_{H S}^{<} p \text { for some } p \in \tau\right\} \tag{1}
\end{equation*}
$$

(Recall that, by the convention in 1.3., (1) means that $I(\tau)$ is generated by the set of seminorms on right hand side.)

If $E$ is a Hilbert space, then $I(\tau)$-topology coincides with so called S-topology which is proven to be of importance in studying characteristic functionals. In fact, there is a complete analogy between the role of S-topology in Hilbert spaces and the role of $I(\tau)$-topology in dual spaces.
2. Notation. Let $(H,\|\cdot\|)$ be a separable Hilbert space. Let $\tau, \tau_{H}$, $\tau_{N}, S$, be topologies defined by:

$$
\begin{aligned}
\tau= & \{\|\cdot\|\} \\
\tau_{H}= & \{p: p(x)=\|A x\|, A \text { is a Hilbert-Schmidt operator }\} \\
\tau_{N}= & \left\{p: p(x)=\langle A x, x\rangle^{\frac{1}{2}}, A \text { is a positive nuclear operator }\right\} \\
S= & \left\{p: p(x)=\langle A x, x\rangle^{\frac{1}{2}}, A\right. \text { is a positive compact, trace } \\
& \text { class operator }\}
\end{aligned}
$$

Operators $A$ in the definition of S-topology are usually called S-operators.
3. Lemma. $\tau_{H}=\tau_{N}=S=I(\tau)$

Proof. Let us first show $\tau_{H}=\tau_{N}$. Let $A$ be a Hilbert-Schmidt operator. Let $p(x)=\|A x\|$. Then $p^{2}(x)=\langle A x, A x\rangle=\left\langle A^{*} A x, x\right\rangle$, and $A^{*} A$ is nuclear by Theorem 1.2.1. Conversely, let $A$ be a positive nuclear operator. Again by Theorem 1.2.1., $A^{\frac{1}{2}}$ is positive Hilbert-Schmidt operator; so $p(x)=\langle A x, x\rangle=\left\|A^{\frac{1}{2}} x\right\|$
${ }^{\tau} N=S$ follows from Theorem 1.2.1.(ii)
To show $I(\tau)=\tau_{H}$, let $P_{A}(x)=\|A x\|$, where $A$ is a Hilbert-Schmidt operator. Then $P_{A}(x) \leq\|A\| \cdot\|x\|$ and $\sum P_{A}^{2}\left(e_{i}\right)=\sum\left\|A e_{i}\right\|^{2}<\infty$, so $P_{A}<{ }_{H S} \tau$. So, $\tau_{H} \subset T(\tau)$

Conversely, if $P \in I(\tau), P \leqslant \delta \tau$, then $P(x) \leq c \cdot\|x\|$ for some $c>0$. There is a map $A: H_{p} \rightarrow H$ such that $p(x)=\|A x\|$ and $\sum\left\|A e_{i}\right\|^{2}=\sum p^{2}\left(e_{i}\right)<\infty$, so $A$ is a Hilbert-Schmidt operator.

Let $(\Omega, F, P)$ be a probability space. In the space of random variables defined on it introduce the topology by the following neighborhood basis at zero:

$$
\begin{equation*}
U\left(\varepsilon, \eta_{1}\right)=\{X: P(\omega:|X(\omega)| \geq \varepsilon) \leq \eta\} \tag{8}
\end{equation*}
$$

The obtained topological space we shall denote by $L_{0}(\Omega, F, P)$. From (8) it follows that $X_{n} \rightarrow 0$ in $L_{o}$ if and only if $X_{n} \rightarrow 0$ in P-probability. Without difficulties we can prove that $X_{n} \rightarrow 0$ in $L_{0}$ if and only if $E\left(\min \left(|X|^{k}, 1\right)\right) \rightarrow 0$ for every $k>0$, and if and only if $E\left(\frac{|X|}{1+|X|}\right) \rightarrow 0$.

Here $X=Y$ if and only if $X=Y$ a.e. [P], i.e. we are considering equivalence classes.
4. Definitions. Let $(E, \tau)$ be a multi-Hilbertian space,
(i) By a random linear functional we mean a linear mapping $X:(E, \tau) \rightarrow$ $L_{0}(\Omega, F, P)$.
(ii) A random linear functional $X$ is called separable if there exists a countable Hilbertien topology $\tau^{\prime}=\tau$ such that $X \in E^{\prime} \tau^{\prime} \subseteq E^{\prime} \tau$. (iii) A random linear functional $X$ is called regular if for every $x \in E, X(x)(\omega)=\left\langle X_{\omega}, x\right\rangle$ where $X_{\omega} \in E_{\tau}^{\prime}$ for every $\omega$. (iv) Random linear functionals $X$ and $Y$ are said to be equivalent if, for every $x$ in $E, P(X(x)=Y(x))=1$.
(v) We say that $X$ is a version of $Y$ if $P(X(x)=Y(x)$ for all $x \in E)=1$.
5. Theorem [14]. A random linear functional $X$ has a $\tau$-regular separable version if and only if $X$ is $I(\tau)$-continuous (i.e. if the mapping $X:(E, I(\tau)) \rightarrow L_{0}(\Omega, F, P)$ is continuous $)$.
6. Theorem (Ito, [14]--Generalized Bochner's theorem)

Let ( $E, \tau$ ) be a multi-Hilbertien space.
Let $f$ be a complex-valued function defined on $E$ such that
(i) $f$ is positive definite,
(ii) $f(0)=1$,
(iii) $f$ is $I(\tau)$-continuous at 0 .

Then
(iv) $f$ is the characteristic function of a separable probability measure $P$ on $E^{\prime}$.

Conversely, (iv) implies (i), (ii) and (iii).
7. Remark. If $E$ is a Hilbert space, then by Lemma 3, $I(\tau)=S$, so Theorem 8 reduces to a well-known result of Sazonov [26]. On the
other hand, when $E$ is a nuclear space, thus $I(\tau)=\tau$, this result is given in [20].

In the next section, we generalize Lévy's continuity theorem to multiHilbertien spaces.

### 2.2 Continuity theorem

The convergence of characteristic functions on real line implies convergence of corresponding probability measures. In infinitely dimensional spaces we have to impose some conditions for relative compactness.

1. Theorem'[23], Ch. VI). Let $H$ be a Hilbert space.

Let $\left\{P_{k}\right\}$ be a sequence of probability measures on $H$, and $f_{k}$ the corresponding sequence of characteristic functionals.
$\left\{P_{k}\right\}$ is relatively compact if and only if for every $\varepsilon>0$, and for every $k=1,2, \ldots$ :
(i) There exists a S-operator $S_{k}$ such that

$$
\begin{align*}
& 1-\operatorname{Ref}_{k}(x) \leq\left\langle S_{k} x, x\right\rangle+\varepsilon \\
& \sup _{k} \sum_{i}\left\langle S_{k} e_{i}, e_{i}\right\rangle\langle\infty  \tag{ii}\\
& \lim _{N} \sup _{k} \sum_{i=N}^{\infty}\left\langle S_{k} e_{i}, e_{i}\right\rangle=0, \tag{iii}
\end{align*}
$$

where $\left\{e_{i}\right\}$ is an orthonormal basis in $H$.
2. Example. (Sazonov, [26]) It is well known that, on real line, a convergent sequence of characteristic functions is equicontinuous. This example shows that it is not true in a Hilbert space. However, as shown in Meyer's paper [20], it remains true in a nuclear space.

Let

$$
t_{i j}=\left\{\begin{array}{ll}
1 / i^{2} j^{2} & \text { if } i \neq j \\
1 / i^{2} & \text { if } i=j
\end{array} \quad i, j=1,2, \ldots\right.
$$

Then

$$
\begin{equation*}
\sum_{j=1}^{\infty} t_{i j}<c<\infty \tag{1}
\end{equation*}
$$

$$
\begin{align*}
& \lim _{N} \sup _{i} \sum_{j=N}^{\infty} t_{i j}=0  \tag{2}\\
& \sum_{j=1}^{\infty} \sup _{i} t_{i j}=\infty \tag{3}
\end{align*}
$$

Define operators $S_{k}, k=1,2, \ldots$ by $s_{k} e_{k}=t_{k j} e_{j}$, where $e_{j}$, $j=1,2, \ldots$ is an orthonormal basis in an infinitely dimensional Hilbert space $H$. Let $P_{k}$ be probability measures on $H$ with characteristic functions $f_{k}(x)=\exp \left(-\frac{1}{2}\left\langle S_{k} x, x\right\rangle\right), k=1,2, \ldots$ By Theorem 1, (2) and (3) imply weak convergence of $f_{k}$.

Suppose now that $\left\{f_{k}\right\}$ is equicontinuous at 0 in $S(=I(\tau))$-topology. Then for every $\varepsilon>0$ there is an $S$-operator $S_{\varepsilon}$ such that $\left\langle S_{\varepsilon} x, x\right\rangle<1$ implies $1-f_{k}(x)<\varepsilon$, for all $k$. By definition of $f_{k}$ it follows $\left\langle S_{k} x, x\right\rangle<\eta(\varepsilon)$ whenever $\left\langle S_{k} x, x\right\rangle<1$ for all $k, \eta(\varepsilon) \rightarrow 0$. Let $t_{\varepsilon, j}=\left\langle\eta(\varepsilon) S_{\varepsilon} e_{j}, e_{j}\right\rangle$.

For every real $r$, if $\left\langle\eta(\varepsilon) S_{\varepsilon} r e_{j}, r e_{j}>=r^{2} t_{\varepsilon, j}<\eta(\varepsilon)\right.$, then $<S_{\varepsilon} r e_{j}, r e_{j}>1$, but then $r^{2} t_{\varepsilon, j}<\eta(\varepsilon)$ implies $r^{2} t_{k j}<\eta(\varepsilon)$ for all $k$. Since $r$ is arbitrary, it follows $t_{k j}<t_{\varepsilon, j}$ for all $k$. Since $S_{\varepsilon}$ is an S-operator we have $\sum_{j} t_{\varepsilon, j}<\infty$, which then implies $\sum \sup _{k} t_{k, j}<\infty$, and this contradicts (3). So, $f_{k}$ is not $I(\tau)-$ equicontinuous.
3. Example. This example is to show that in the infinitely dimensional Hilbert space, the convergence of characteristic functionals does not imply convergence of corresponding measures.

Let $t_{i j}$ be any set of real numbers satisfying
(4)

$$
t_{i j} \rightarrow u_{j} \quad \text { as } i \rightarrow \infty, j=1,2, \ldots
$$

$$
\begin{equation*}
0<t_{i j} \leq 1 \tag{5}
\end{equation*}
$$

$$
\lim _{N \rightarrow \infty} \sup _{i} \sum_{j=N}^{\infty} t_{i j}>0
$$

For an orthonormal sequence $\left\{e_{i}\right\}$ in $H$, define the operators $S_{k}$ and $S$ by $S_{k} e_{j}=t_{k j} e_{j} ; S e_{j}=u_{j}, j, k=1,2, \ldots$

By (4) we have $\left\langle S_{k} x, x\right\rangle \rightarrow\langle S x, x\rangle$ for every $x$ in $H$. Define $f_{k}(x)=\exp \left\{-\frac{1}{2}\left\langle S_{k} x, x>\right\}\right.$ and $f(x)=\exp \left\{-\frac{1}{2}<S x, x>\right\}$. Then clearly $f_{k}(x) \rightarrow f(x)$ for every $x$. But by (6), the corresponding sequence of measures is not weakly convergent.

Until further notice, E will denote a milti-Hilbertien space, with topology $\tau$.
4. Theorem.

Let $M_{k}$ be a sequence of separable probability measures on $E^{\prime}$, let $f_{k}$ be the corresponding sequence of characteristic functionals. Assume the following:
(i) There is a function $f, I(\tau)$-continuous at 0 , such that

$$
\begin{equation*}
f_{k}(x) \rightarrow f(x), \text { for every } x \text { in } E . \tag{7}
\end{equation*}
$$

(ii) For every $\varepsilon>0$ there is a sequence of $I(\tau)$-seminorms $p_{k}$ and a r-seminorm $q$ such that $p_{k} \ll q$ for every $k$ and
(8)

$$
1-\operatorname{Re} f_{k}(x) \leq \varepsilon+p_{k}^{2}(x)
$$

$$
\begin{equation*}
\sup _{k} \sum_{i} p_{k}^{2}\left(e_{i}\right)<\infty \tag{9}
\end{equation*}
$$

$$
\begin{equation*}
\lim _{N} \sup _{k} \sum_{i=N}^{\infty} p_{k}^{2}\left(e_{i}\right)=0 \tag{10}
\end{equation*}
$$

where $\left\{e_{i}\right\}$ is an orthonormal basis in $E_{q}$.
Then there is a separable probability measure $M$ on $E^{\prime}$ such that $M_{k} \Rightarrow M$.

Conversely, if $M_{k} \Rightarrow M$, where $M$ is a separable probability measure, then (7)-(10) hold.

Proof. Assume (7)-(10). By Theorem 1, the set of measures induced by $M_{k}$ in $E_{q}$ is relatively compact. So there exists a $q$-compact (thus $\tau$-compact) set $K$ such that $M_{k}\left(K^{c}\right) \leq 1-\varepsilon$.

So, $M_{k}$ is a tight sequence for which $f_{k}(x) \rightarrow f(x)$ pointwise.
By Theorem 1.5, there is a separable probability measure $M$ such that
$f$ is the characteristic functional of $M$; thus $M_{k} \Rightarrow M$.
Conversely, let $M_{k} \Rightarrow M$. Then there is a compact set $K \subseteq E_{\tau}^{\prime}$, such
that $M_{k}\left(K^{C}\right)<\varepsilon / 2$ for every $k=1,2, \ldots$.
For every $x$ and every $k$ we have
$1-\operatorname{Re} f_{k}(x)=\int(1-\cos \langle F, x\rangle) d M_{k}(F)$

$$
\begin{aligned}
& \leq \int_{K}(1-\cos \langle F, x\rangle) d M_{k}(F)+\varepsilon \\
& \leq \frac{1}{2} \int_{K}\langle F, x\rangle^{2} d M_{k}+\varepsilon
\end{aligned}
$$

$K$ is compact in $E_{\tau^{\prime}}^{\prime}$; so there is a $q$ such that $K$ is compact in $E_{q}^{\prime}$.

Let $p_{k}^{2}(x)=\frac{1}{2} \int_{K}\langle F, x\rangle^{2} d M_{K}(F)$.
If $\left\{e_{i}\right\}$ is an orthonormal basis in $E_{q}$ then

$$
\sum_{i} p_{k}^{2}\left(e_{i}\right)=\frac{1}{2} \int_{K} \bar{q}^{-2}(F) d M_{K}(F)<\frac{1}{2} \sup _{F \in K} \bar{q}^{-2}(F)<\infty
$$

because $\bar{q}$ is bounded on $K$, being a continuous function. Also we have

$$
p_{k}^{2}(x) \leq q^{2}(x) \cdot \frac{1}{2} \int_{k} q^{-2}(F) d M_{k}(F) \leq C \cdot q^{2}(x) \text {, }
$$

so $p_{k}<{ }_{H S} q$, for every $k$. Thus, (8) and (9) are proved.
To prove (10) note that

$$
\begin{aligned}
\sum_{i=N}^{\infty} p_{k}^{2}\left(e_{i}\right) & =\sum_{N}^{\infty} \int_{K}\left\langle F, e_{i}\right\rangle^{2} d M_{k}(F) \\
& \leq \sup _{F \in K} \sum_{N}^{\infty}\left\langle F, e_{i}\right\rangle^{2} \rightarrow 0, \text { as } N \rightarrow \infty,
\end{aligned}
$$

by compactness of K.
5. Corollary. Let $f_{k}, k=1,2, \ldots$, be a sequence of characteristic functionals of separable probability measures $M_{k}, k=1,2, \ldots$ on $E^{\prime}$, which is $I(\tau)$-equicontinuous at 0 , i.e., there is a seminorm $p$ in $I(\tau)$ such that for every $\varepsilon>0$ there is a $\delta>0$ so that for all $k=1,2, \ldots$ we have:

$$
\begin{equation*}
p(x) \leq \delta \Rightarrow 1-\operatorname{Ref}(x) \leq \varepsilon . \tag{11}
\end{equation*}
$$

Then $\left\{M_{k}\right\}$ is a relatively compact sequence.
Proof. Assume (11). Then conditions (7)-(10) are satisfied for $p_{k}=\frac{1}{\delta} \cdot p, k=1,2, \ldots$, and $q$ such that $p<q$. Then from the proof of Theorem 4 it follows that $\left\{M_{k}\right\}$ is relatively compact.
6. Remark. Theorem 4 and Corollary 5 remain true for any sequence of finite separable measures $M_{n}$ such that $\sup _{n} M_{n}\left(E_{\tau^{\prime}}^{\prime}\right)<\infty$.

### 2.3. Infinitely divisible probability measures

Let ( $E, \tau$ ) be a multi-Hilbertian space.

We have defined infinitely divisible probability measures in Chapter 1. By Theorem 1.5., a separable probability measure on $E^{\prime}$ is uniquely determined by its characteristic function, so we have

1. Theorem. Let $M$ be a separable probability measure on $E^{\prime}$, with the characteristic function $f . M$ is an ID measure if for every $n$ there is a characteristic function $f_{n}$ such that

$$
\begin{equation*}
f_{n}^{n}(x)=f(x), \tag{12}
\end{equation*}
$$

for every $x$ in $E$.

The following theorem can be proved in the same way as for the real random variables.
2. Theorem. (i) ID characteristic functional never vanishes.
(ii) If $f$ is an ID characteristic functional, then $f^{s}$ is ID characteristic functional, for every s > 0 .
(iii) The sum of finite number of ID random variables is ID.
(iv) Weak limit of ID measures is an ID measure; moreover, if only $f_{k} \rightarrow f$ pointwise, and $f_{k}$ are ID characteristic functionals, then $f$ is an ID characteristic functional.
3. Examples. (i) Let

$$
\begin{equation*}
f(x)=\exp C\left(\int e^{i<F, x\rangle} d M(F)-1\right), \tag{13}
\end{equation*}
$$

where $M$ is a separable probability measure on $E^{\prime}$, and $C$ any real number.

The function defined by (13) is positive definite, $I(\tau)$-continuous and
$f(0)=1$. So, by Theorem 1.5 it is the characteristic functional of a separable probability measure on E'-Poisson measure. It satisfies (12), so it is ID measure.
(ii) Let

$$
\begin{equation*}
\left.f(x)=\exp (i<F, x\rangle-\frac{1}{2} p^{2}(x)\right) \tag{14}
\end{equation*}
$$

where $F$ is a fixed element in $E^{\prime}, P$ is a $I(\tau)$-seminorm. By Bochner's theorem there is a unique separable probability measure on $E^{\prime}$ determined by (14)--Gaussian measure. It is clearly an ID measure. In Chapter 1, the construction of a Gaussian measure, starting from cylinder measures is given.
4. Theorem. A function $g$ is the logarithm of an ID characteristic functional $f$ (of a separable probability measure on $E^{\prime}$ ) if and only if

$$
\begin{equation*}
g(0)=0 \tag{15}
\end{equation*}
$$

$$
\begin{equation*}
g \text { is } I(\tau) \text {-continuous at } 0, \text { and } \tag{16}
\end{equation*}
$$

for all finite sets of complex numbers
$A_{I}=\left\{a_{i}\right\}_{i \in I}$ such that $\sum_{i \in I} a_{i}=0$, and $X_{I}=\left\{x_{i}\right\}_{i \in I}$
$x_{i} \in E$, we have

$$
\begin{equation*}
\sum_{i, j \in I} a_{i} \bar{a}_{j} g\left(x_{i}-x_{j}\right) \geq 0 \tag{17}
\end{equation*}
$$

Proof. The assertion follows directly from the Theorem 1.5. In fact, (17) is the consequence of positive definiteness of $f$ and Schoenberg's theorem [19], (16) follows from $I(\tau)$-continuity of $f$, and (15) from $f(0)=1$.
5. Remark. Making an appropriate choice of $A_{I}$ and $X_{I}$ in (17), one can show that for every $x \in E$

$$
\begin{align*}
& \text { Re } g(x) \leq 0  \tag{18}\\
& g(-x)=\overline{g(x)} . \tag{19}
\end{align*}
$$

6. Theorem. The class of ID characteristic functionals coincide with the class of Poisson characteristic functionals and their pointwise limits.

Proof. By Example 3 and Theorem 2, Poisson characteristic functionals and their limits are ID.

Conversely, let $f$ be an ID characteristic functional.
Then for all $x$ in $E$,

$$
n\left(f^{1 / n}(x)-1\right) \rightarrow \log f(x) \quad(n \rightarrow \infty)
$$

so, $f(x)=\lim _{n} \exp \left(n\left(f^{1 / n}(x)-1\right)\right)$. By Theorem $1, f^{1 / n}$ is a characteristic functional of some separable probability measure $M_{n}$, so

$$
f^{l / n}(x)=\int e^{i\langle F, x\rangle} d M_{n}(F)
$$

and

$$
f(x)=\lim _{n} \exp n \int\left(e^{i\langle F, x\rangle}-1\right) d M_{n}(F),
$$

which is a limit of Poisson characteristic functionals.

Now we give some elementary inequalities that will be used later.
7. Lemma. For any real $t, s$ :
(i) $\left|e^{i t}-1\right| \leq 2 \min \left(t^{2}, 1\right) \leq 2 \min (t, 1)$
(ii) $\left|e^{i t}-1-i t\right| \leq t^{2} / 2$
(iii) $1-\sin t / t \geq c \cdot \min \left(t^{2}, 1\right)$
(iv) $\left|\left(e^{i s}-i s\right)-\left(e^{i t}-i t\right)\right| \leq|s-t| \cdot \frac{|s|+|t|}{2}$

Proof: elementary.
3. Notation. For $q \in \tau$ define the following quantity (finite or not)

$$
\begin{equation*}
\tilde{q}(F)=\sup _{q(x) \leq 1}|<F, x>| \tag{20}
\end{equation*}
$$

Let $M$ be a positive separable measure on $E^{\prime}-\{0\}$ such that

$$
\begin{equation*}
E_{E^{\prime}-\{0\}}^{\int} \min \left(\sim^{2}(F), 1\right) d M(F) .<\infty . \tag{21}
\end{equation*}
$$

For a $G$ in $E^{\prime}$ and $r \in I(\tau)$ define

$$
\begin{align*}
g(x) & =g[G, r, M, q](x)  \tag{22}\\
& \left.=i\langle G, x\rangle-r^{2}(x)+\int\left(e^{i<F, x\rangle}-1-i<F, x\right\rangle \cdot 1[\tilde{q}(F) \leq 1]\right) d M(F)
\end{align*}
$$

Function defined by (22) will be referred to as $g[G, r, M, q]$ or only $g$. We will also use the following notation:

$$
\begin{equation*}
h(F, x)=e^{i<F, x\rangle}-1-i\langle F, x\rangle \cdot 1[\tilde{q}(F) \leq 1] \tag{23}
\end{equation*}
$$

9. Lemma. $g$ is $I(\tau)$-continuous at 0 .

Proof: $G \in E^{\prime}$ implies $G \in E_{p}^{\prime}$ for some $p \in \tau$.
Let $\left\{e_{j}\right\}$ be an orthonormal basis in $E_{p}$.
Then $\left[\left\langle G, e_{i}\right\rangle^{2}=\tilde{p}(G)\right.$, thus $|\langle G, x\rangle|$ is a $I(\tau)$-continuous seminorm, so $i\langle G, x\rangle$ is $I(\tau)$-continuous.

Since $r \in I(\tau)$, it remains to show the $I(\tau)$-continuity of the integral in (22). We have:

$$
\begin{equation*}
\int h(F, x) d M(F)=\underset{\sim}{\tilde{q}(F)>1} \underset{f}{f}\left(e^{i<F, x\rangle}-1\right) d M(F)+\underset{\tilde{q}(F) \leq 1}{\int}\left(e^{i\langle F, x\rangle}-1-i\langle F, x\rangle\right) d M(F) \tag{24}
\end{equation*}
$$

Let us show that both terms above are $I(\tau)$-continuous.
By Lemma 7(i),

$$
\begin{equation*}
\tilde{\tilde{q}}_{\int}^{\int}(F)>1 \mathrm{e}\left(e^{i\langle F, x\rangle}-1\right) d M(F) \mid \leq \underset{\tilde{q}(F)>1}{\int} \min \left(\langle F, x\rangle^{2}, 1\right) d M(F) \tag{25}
\end{equation*}
$$

Restricted to the set where $\tilde{\mathfrak{q}}(F)>1, M$ is a positive finite measure;
let $M(\tilde{q}(F)>1)=m$. By separability of $M$, there is a topology $\sigma$ determined by an increasing family of seminorms $q_{n}$ such that $E_{q_{n}^{\prime}}^{\prime} \uparrow E_{\sigma}^{\prime}$. So, for given $\varepsilon$ we can find $n$ such that

$$
M\left(F \notin E^{\prime} a_{n}\right) \leq \varepsilon / 2
$$

i.e., if $\left\{e_{k}\right\}$ is $a_{n}$-orthonormal basis:

$$
\left.M\left(\sum<F, e_{k}\right\rangle^{2}=\infty\right) \leq \varepsilon / 2 .
$$

Let $r=r(\varepsilon)$ be a real number such that

$$
\begin{equation*}
M\left(\sum<F, e_{k}>^{2} \geq r\right) \leq \varepsilon / 2 . \tag{26}
\end{equation*}
$$

Then, with $F_{k}=\left\langle F, e_{k}\right\rangle$, we have
(27) $\underset{\tilde{q}(F)>1}{\int \min }\left(\langle F, x\rangle^{2}, 1\right) d M \leq \underset{\widetilde{q}(F)>1}{\int}\langle F, x\rangle^{2} \cdot 1\left[\sum F_{k}^{2}<r_{i}^{-} d M(F)\right.$

$$
+\int 1\left[\sum F_{k}^{2} \geq r\right] d M(F)
$$

Let $p^{2}(\varepsilon, x)=p^{2}(x)=\underset{\widetilde{q}(F)>1}{\int}\left\langle F, x>^{2} \cdot 1\left[\sum F_{k}^{2}<r\right] d M(F)\right.$
Then $\sum p^{2}\left(e_{k}\right)<r \cdot m$; so $p$ is an $I(\tau)$-seminorm ; by (25) and (27) we have

$$
\left|\underset{\widetilde{q}(F)>1}{\int}\left(e^{i<F, x>}-1\right) d M(F)\right| \leq p^{2}(x)+\varepsilon,
$$

so the first term in (24) is $I(\tau)$-continuous.
By Lemma 7.(ii):
(28) $\left.\mid \underset{\widetilde{q}(F)<1}{\int} e^{i<F, x>}-1-i<F, x>\right) d M(F)\left|\leq \frac{1}{2} \underset{\widetilde{q}(F)<1}{\int}<F, x\right\rangle^{2} d M(F)$.

Now, by definition of $\mathfrak{q}$, it follows

$$
|<F, x\rangle \mid \leq \tilde{q}(F) q(x) \text {, }
$$

so by the assumption (21), the right hand side of (28) is finite. Moreover, on the set where $\tilde{q}(F)<1$, we have $F \in E_{q}^{\prime}$, so the expression on the
lefthand side of (28) is $I(\tau)$-continuous.
10. Corollary. $g$ is logarithm of an ID characteristic functional of a separable probability measure on $E^{\prime}$.

Proof. We shall use Theorem 4. Condition (15) is satisfied; (16) follows from Lemma 9 and (17) can be easily checked.
11. Lemma. For given $q \in \tau$, there is $1-1$ correspondence between functions $g$ defined by (22) and triplets [G, r, M].

Proof: Let $t \in R$. We have

$$
\begin{aligned}
\left|\operatorname{Re} \frac{e^{i\langle F, t x\rangle}-1}{t^{2}}\right| \leq\left|\frac{e^{i<F, t x\rangle-1}}{t^{2}}\right| & \leq \min \left(\langle F, x\rangle^{2}, 2 / t^{2}\right) \\
& \leq 2 q^{2}(x) \cdot \min \left(\tilde{q}^{2}(F), 1\right),
\end{aligned}
$$

for a large enough $t$.

So, by dominated convergence theorem,

$$
t^{-2} \cdot \int_{\tilde{q}(F)>1}\left(e^{i<F, t x\rangle}-1\right) \cdot d M(F) \rightarrow 0 \text { as } t \rightarrow \infty
$$

Similarly,

$$
t^{-2} \cdot \int_{\tilde{q}(F) \leq 1}\left(e^{i<F, t x>}-1-i<F, t x>\right) \quad d M(F) \rightarrow 0 \text { as } t \rightarrow \infty
$$

So, we have

$$
r^{2}(x)=\lim _{t \rightarrow \infty}-\operatorname{Reg}(t x) \cdot t^{-2} \quad(\geq 0 \text { by }(18)),
$$

which shows that $r$ is uniquely determined by $g$.

Let now $x, y \in E$. The following formula hoids:

$$
\begin{equation*}
\int(1-\cos \langle F, y\rangle) e^{i\langle F, x\rangle} d M(F)=r^{2}(y)+g(x)-\frac{g(x+y)+g(x-y)}{2} . \tag{29}
\end{equation*}
$$

We shall show that $M$ is uniquely determined by (29).
For $x=0$ and $y$ being fixed, we have from (29) that the measure $N$ defined by

$$
\begin{equation*}
d N=(1-\cos \langle F, y\rangle) d M \tag{30}
\end{equation*}
$$

is a positive finite measure on $E^{\prime}$ (more precisely, its extension to $E^{\prime}$ ) Since $N$ has a $I(\tau)$-continuous characteristic functional, it follows that $N$ is a separable finite measure, uniquely determined by (30). So, $M$ is uniquely determined by $N$ on any set on which $1-\cos \langle F, y\rangle$ is bounded away from 0 . Since $y$ runs through $E$, by regularity we conclude that $M$ is uniquely determined. Finally, $G$ is determined by $M, q$ and $r$, from (22).

The following two lemmas are proved in [6], in a different context:
12. Lemma. Let $f$ be an ID characteristic functional.

Let $f_{n}=f^{1 / n}$. Suppose that, for some $\varepsilon>0$ and some $p \in I(\tau)$, $p<\mathrm{HS}$, we have:

$$
\begin{equation*}
p(x) \leq 1 \Rightarrow 1-f(x) \leq \varepsilon . \tag{31}
\end{equation*}
$$

Then
(32) $n\left(1-\operatorname{Re} f_{n}(x)\right) \leq 8 \varepsilon(1+p(x))$, for all $x$, and
(33) $\int \min \left(q^{2}(F), 1\right) d U_{n} \leq 48 \varepsilon$, for ait $n$,
where $U_{n}$ is the measure that corresponds to $n f_{n}$.
13. Lemma. Let $Q=\{F: q(F) \leq 1\}-\{0\}$. Let $K(F, x)=$
$=\tilde{q}^{-2}(F)\left(e^{i<F, x\rangle}-1-i\langle F, x\rangle\right)$. Let $M_{n}$ be a sequence of positive
finite measures on $Q$, such that $M_{n}(Q) \leq 1$ and $\int K(F, x) d M_{n}(F)$
converges pointwise to a function $w(x)$. Then there is a positive separable measure $M_{0}$ on $Q$, such that $M_{0}(Q) \leq 1$ and an $I(\tau)$-seminorm $r$, such that $w(x)=-r^{2}(x)+\int K(F, x) d M_{0}(x)$.
14. Theorem. Let $g$ be the logarithm of an ID characteristic functional of a separable probability measure on $E^{\prime}$. Suppose $q \in \tau$ such that there is a $p \in I(\tau), p \underset{H S}{<} q$, and

$$
\begin{equation*}
p(x) \leq 1 \Rightarrow|g(x)| \leq \varepsilon \tag{34}
\end{equation*}
$$

for some $\varepsilon, 0<\varepsilon<\frac{1}{4}$.
Then there is a triplet $[G, r, M]$ satisfying conditions of Notation 8, such that $g=g[G, r, M, q]$. This representation is unique.

Proof. Let $g=\log f$, where $f$ is an ID characteristic functional of a separable probability measure on $E^{\prime}$. Let $f_{n}=n f^{1 / n}$. Let $Q$ be as in Lemma 13. We have $\log f=g=\lim _{n} n\left(f^{1 / n}-1\right)$. Let $P_{n}$ be the measure corresponding to $f_{n}$. Let $U_{n}$ be $P_{n}$ restricted to $Q$ : let $V_{n}$ be $P_{n}$ restricted to $Q^{C}-\{0\}$. Denote by $\hat{U}_{n}$ and $\hat{V}_{n}$ the corresponding characteristic functionals. Let $f(x)=1-Z$. Then, as $n \rightarrow \infty$ we have

$$
\begin{aligned}
f_{n}(0)-f_{n}(x) & =n\left(1-f^{1 / n}(x)\right)=n\left(1-(1-z)^{1 / n}\right) \\
& \doteq n\left(1-\left(1-\frac{z}{n}\right)\right) \doteq-z .
\end{aligned}
$$

so we conclude that $\left\{f_{n}\right\}$ is equicontinuous at 0 . Now we have:

$$
f_{n}(0)-\operatorname{Re} f_{n}(x)=\left(\hat{U}_{n}(0)-\operatorname{Re} \hat{U}_{n}(x)\right)+\left(\hat{V}_{n}(0)-\operatorname{Re} V_{n}(x)\right)
$$

and by $\hat{U}_{n}(0)-\operatorname{Re} \hat{U}_{n}(x) \geq 0$, we conclude that

$$
f_{n}(0)-\operatorname{Re} f_{n}(x)<\varepsilon \Rightarrow \hat{V}_{n}(0)-\operatorname{Re} \hat{V}_{n}(x)<\varepsilon
$$

i.e., $\hat{V}_{n}$ is equicontinuous at 0 .

By (34) and Lemma 12 we have

$$
\begin{equation*}
V_{n}\left(E^{\prime}\right)=P_{n}\left(Q^{c}-\{0\}\right) \leq 48 \varepsilon, \tag{35}
\end{equation*}
$$

so by Corollary 2.5. and Remark 2.6. it follows that $V_{n}$ is relatively compact. Let $V$ be a measure such that $V_{n}{ }^{\prime} \Rightarrow V$ for some subsequence $n^{\prime}$. Then $\hat{V}_{n}{ }^{\prime}-\hat{V}_{n}(0) \rightarrow \hat{V}-\hat{V}_{n}(0)$ pointwise, By (34) and Lemma 12,

$$
\begin{equation*}
\int \tilde{q}^{2} d U_{n} \leq 48 \varepsilon \tag{36}
\end{equation*}
$$

By inequality 2.7.(ii) we have, for every $x$ in $E$ :

$$
\begin{aligned}
& \left.\mid \int\left(e^{i<F, x\rangle}-1-i<F, x\right\rangle\right) d U_{n^{\prime}}(F) \mid \\
& \quad \leq \int \frac{\langle F, x\rangle^{2}}{2} d U_{n^{\prime}}(F) \\
& \quad \leq 12 \varepsilon q^{2}(x) .
\end{aligned}
$$

Now let $w_{n}(x)=\int\left(e^{i\langle F, x\rangle}-1-i\langle F, x\rangle\right) d U_{n}{ }^{\prime}$.
The above inequality shows that $w_{n}(x)$ is bounded for every $x$; therefore, there is a function $w(x)$ such that for every $x$, $w_{n \prime \prime}(x) \rightarrow w(x)$, for some subsequence $n^{\prime \prime}$.

We have obtained so far:

$$
\begin{aligned}
g(x) & =\lim _{n} n\left(f_{n}(x)-1\right) \\
& =\lim _{n^{\prime \prime}}\left(\left(\hat{V}_{n^{\prime \prime}}(x)-\hat{V}_{n \prime \prime}(0)\right)+\left(\hat{U}_{n^{\prime \prime}}(x)-\hat{U}_{n^{\prime \prime}}(0)\right)\right) \\
& =(\hat{V}(x)-\hat{V}(0))+w(x)+i \lim _{n^{\prime \prime}} f<F, x>d U_{n^{\prime \prime}}(F) \\
& =\hat{V}(x)-\hat{V}(0)+w(x)+i\left\langle G_{0}, x\right\rangle .
\end{aligned}
$$

By Lemma 13, we have that there is a $r \in I(\tau)$, a separable measure $M_{0}$ on $Q$ such that

$$
\left.w(x)=-r^{2}(x)+\int\left(e^{i<F, x\rangle}-1-i<F, x\right\rangle\right) \cdot \tilde{q}^{-2}(F) d M_{0}
$$

To conclude the proof, define $G$ by

$$
\langle G, x\rangle=\left\langle G_{0}, x\right\rangle+\underset{\widetilde{G}(F) \leq 1}{\int}\langle F, x\rangle d V(F)
$$

and define measure $M$ on $E^{\prime}-\{0\}$ to be $M_{0}$ on $Q$ and $V$ on $Q^{C}-\{0\}$. Clearly, the condition (21) is satisfied and $g=g[G, r, M, q]$.
Uniqueness is proved in Lemma 11.

## 3. CONVERGENCE OF ID LAWS AND HOMOGENEOUS PROCESSES WITH INDEPENDENT INCREMENTS

### 3.1. Convergence of ID laws

1. Definition. We say that a sequence of $I(\tau)$-seminorms $\left\{p_{n}\right\}$ is compact if there exists a $\tau$-seminorm $t$ such that $p_{n}<t$ for all $n$ and, for an orthonormal basis $\left\{e_{j}\right\}$ in $E_{t}$, we have:

$$
\begin{equation*}
\sup _{n} \sum_{i} p_{n}^{2}\left(e_{i}\right)<\infty \tag{1}
\end{equation*}
$$

$$
\begin{equation*}
\lim _{N} \sup _{n} \sum_{i=N}^{\infty} p_{n}^{2}\left(e_{i}\right)=0 \tag{2}
\end{equation*}
$$

2. Theorem. Let $P_{n}$ and $P$ be separable ID probability measures on $E$ ' with the characteristic functionals $g_{n}=g_{n}\left[G_{n}, r_{n}, M_{n}, q\right]$ and $g=g[G, r, M, q]$. Then $P_{n} \Rightarrow P$ if and only if

$$
\begin{equation*}
\lim _{n} G_{n}=G \tag{3}
\end{equation*}
$$

$$
\begin{equation*}
M_{n} \Rightarrow M \text { on every set }\{F: \tilde{q}(F)>\varepsilon\} \tag{4}
\end{equation*}
$$

(5) The sequence $\left\{t_{n}\right\}$ of $I(\tau)$-seminorms defined by

$$
t_{n}^{2}(x)=r_{n}^{2}(x)+\underset{\tilde{q}(F) \leq 1}{\int}<F, x>^{2} d M_{n}(F)
$$

is compact,
(6) for every $x, \lim _{\varepsilon \rightarrow 0} \operatorname{Tim}_{n \rightarrow \infty} \frac{1}{2} \underset{\tilde{q}(F) \leq \varepsilon}{\int}\left\langle F, x>^{2} d M_{n}(F)+r_{n}^{2}(x)=r^{2}(x)\right.$.

Proof. Assume the conditions above. Then for every $x$ in $E$

$$
\begin{align*}
g_{n}(x) & \left.=i\left\langle G_{n}, x\right\rangle+\underset{\tilde{q}(F) \geq \varepsilon}{\dot{\tilde{q}}\left(e^{i<F, x\rangle}-1-i<F, x\right\rangle} \cdot 1[\tilde{q}(F) \leq 1]\right) d M_{n}(F)  \tag{7}\\
& +\underset{\tilde{q}(F) \leq \varepsilon}{\int}\left(e^{i<F, x\rangle}-1-i<F, x>\right) d M_{n}(F)-r_{n}^{2}(x) .
\end{align*}
$$

Now we let $n \rightarrow \infty$ to obtain that first term in (7) converges to $i<G, x>$
and the second term to the corresponding expression with $M$. Now we take care of the remainder. Use inequality 2.3.7.(ii) and note that if $|a|<|b|$ then $|a-c|<|-b-c|$, for $a, b, c$ complex numbers. So we have
(8)

$$
\begin{aligned}
& \left.\int\left(e^{i<F, x\rangle}-1-i<F, x\right\rangle\right) d M_{n}(F)-r_{n}^{2}(x)+r^{2}(x) \mid \leq \\
& \tilde{q}(F) \leq \varepsilon \\
& \left.\leq\left|-\frac{1}{2} \underset{\tilde{q}(F) \leq \varepsilon}{\int}<F, x\right\rangle^{2} d M_{n}(x)-r_{n}^{2}(x)+r^{2}(x) \right\rvert\, .
\end{aligned}
$$

Letting now $\varepsilon \rightarrow 0$, and using (6) we have that for every $x$ in $E$, $g_{n}(x) \rightarrow g(x)$, so then characteristic functionals $f_{n}$ converge pointwise to the characteristic functional $f$ of the measure $M$.

The condition (5), together with inequalities in the proof of Lemma 2.9., provides condition (ii) of Theorem 2.4. So, by Theorem 2.4., we conclude that $P_{n}=P$.

Conversely, let $P_{n} \Rightarrow P$. Using Theorem 2.4. and the proof of Theorem 5.5. of [23] we obtain conditions (3)-(6).

### 3.2. Homogeneous processes with independent increments

Let $T$ be a finite or infinite interval on the real line, starting at 0 .

1. Definition. (i) A process with independent increments is a family of random variables $\{X(t)\}_{t \in T}$, defined on a probability space ( $\Omega, F, P$ ), such that for every $t_{1}, \ldots, t_{n}\left(0 \leq t_{1}<\ldots<t_{n}\right)$, the random variables

$$
x\left(t_{0}\right), x\left(t_{1}\right)-x\left(t_{0}\right), \ldots, x\left(t_{n}\right)-x\left(t_{n-1}\right)
$$

are independent.
(ii) $\{X(t)\}_{t \in T}$ is said to be a homogeneous process
if the distribution of random variable $X(t)-X(s)(s<t)$ depends on
t - s only.

A connection between homogeneous processes with independent increments and ID laws is immediate. Let $t \in T$ be fixed. By Definition 1 , $x\left(\frac{k t}{n}\right)-x\left(\frac{(k-1) t}{n}\right)$ are fid random variables $(k=1, \ldots, n)$, and $x(t)-x(0)=\sum_{k=1}^{n} x\left(\frac{k t}{n}\right)-x\left(\frac{(k-1) t}{n}\right)$ is an ID random variable. Therefore, every increment of a homogeneous process with independent increments is an ID random variable. Let $P_{t}$ be the distribution of $X(t)-X(0)$, and $Q$ be the distribution of $X(0)$. It is easy to see that $P_{t+s}=P_{t}^{*} P_{s}$ and $P_{o}=\delta_{0}$ (Dirac distribution), so $\left\{P_{t}\right\}_{t \in T}$ is a semigroup of convolutions. So, the distribution of $X_{t}$ is obtained as $P_{t}{ }^{*} Q$; with some additional work it can be shown ([1]) that all finitely dimensional distributions of the process $X_{t}$ are determined by the semigroup $\left\{P_{t}\right\}$ and $Q$.

Let now $\left\{X_{t}\right\}$ be a $E^{\prime}$-valued process, and let $f_{t}$ be the characteristic functional of $X_{t}$. Until further notice assume that $\left\{X_{t}\right\}$ is a homogeneous process with independent increments.

By independence we have:

$$
\begin{equation*}
f_{t+s}(x)=f_{t}(x) \cdot f_{s}(x) \tag{1}
\end{equation*}
$$

In a particular case when $f_{t}(x)$ is, for every $x$, continuous at $t=0$, we have an especially simple relation.

1. Theorem. Suppose that for every $x$ in $E, f_{t}(x)$ is continuous at $t=0$. Then

$$
\begin{equation*}
f_{t}(x)=\left[f_{1}(x)\right]^{t} \tag{2}
\end{equation*}
$$

Proof: Let $s \rightarrow 0$ for fixed $t$ in (1) to show that $f$ is continuous
for every $t$. The only continuous solutions of (1) are functions of the form

$$
\begin{equation*}
f_{t}(x)=\exp t g(x) \tag{3}
\end{equation*}
$$

for some g. This proves (2).
Suppose now that $f_{1}$ is the characteristic functional of a separable probability measure on $E^{\prime}$. Then (2) completely determines the process. If $X(0)=0$ with probability 1 , then for every $t, X_{t}$ is an ID random variable; by (3) and results in Chapter 2, we know the form of $f_{t}(x)$, and by Theorem 1.2. we have necessary and sufficient conditions for weak convergence.

Note that $f_{1}$ is an ID characteristic functional.

### 3.3. An Example

In this section we present the solution to a problem in real line, which arises in a stochastical model of neuronal activity. This is an improvement and a generalization of results of Kallianpur [16] and Tuckwell [29]. We also discuss a possiblity of a generalization to infinitely dimensional spaces.

Let us first recall some facts about real valued ID random variables.
Let $X$ be a random variable with the characteristic function $f$, and its logarithm $g$. Suppose that $\operatorname{Var} X<\infty$. Then $g$ is represented in the form

$$
\begin{equation*}
g(x)=i G x+\int\left(e^{i u x}-1-i u x\right) u^{-2} d K(u) \tag{1}
\end{equation*}
$$

where $K$ is distribution function of some finite measure on $R$. The representation (1) is unique and we write $g=g[G, K] . X_{n} \Rightarrow X$ if and only if $K_{n}(u) \rightarrow K(u)$, for every $u \in R$, and $G_{n} \rightarrow G$.

If $X$ is normal with mean $m$ and the standard deviation $S$, then $G=m$ and $K(u)=s^{2} \cdot 1[u>0]$

If $X$ is generalised Poisson, i.e., if

$$
\begin{equation*}
P(X=a+k h)=\frac{e^{-\lambda}}{k!} \lambda^{k}, k=0,1,2, \ldots \tag{2}
\end{equation*}
$$

then $G=a+\lambda h ; k(u)=\lambda h^{2} \cdot 1[t \geq h]$
Let $Y_{k n}(t) \quad k=1,2, \ldots, p_{n}, n=1,2, \ldots$ be independent Poisson processes with parameters $\lambda_{k n}$; let $\varepsilon_{k n}, k=1,2, \ldots, p_{n}$, $n=1,2, \ldots$ be real numbers, $0 \leq t \leq T<\infty$.

Define

$$
\begin{equation*}
N_{n}(t)=\sum_{k=1}^{p_{n}} \varepsilon_{k n} Y_{k n}(t) \tag{3}
\end{equation*}
$$

$$
\begin{equation*}
\bar{N}_{n}(t)=\left(N_{n}(t)-E N_{n}(t)\right) \cdot \sigma_{n}^{-1} \tag{4}
\end{equation*}
$$

where $E N_{n}(t)=t \cdot \sum_{k=1}^{p_{n}} \varepsilon_{k n} \lambda_{k n} ; \sigma_{n}=\left(\sum_{k=1}^{p_{n}} \varepsilon_{k n}^{2} \lambda_{k n}\right)^{\frac{1}{2}}$
Then for $N_{n}(1)$ we have the representation:

$$
G_{n}=\int_{1}^{p_{n}} \lambda_{k n} \varepsilon_{k n}, \quad k_{n}(u)= \begin{cases}\lambda_{k n} \varepsilon_{k n}^{2} & \text { if } u=\varepsilon_{k n}  \tag{5}\\ 0 & \text { otherwise }\end{cases}
$$

where $k_{n}$ is the point mass function;

$$
k_{n}(u)=\sum_{v<u} k_{n}(v)
$$

Similarly, $\bar{N}_{n}(1)$ has the representation:

$$
G_{n}=0 ; k_{n}(u)= \begin{cases}\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \cdot \lambda_{k n} & \text { if } u=\frac{\varepsilon_{k u}}{\sigma_{n}}  \tag{6}\\ 0 & \text { otherwise }\end{cases}
$$

The interest is to investigate the limit behavior of $N_{n}(t)$ and $\bar{N}_{n}(t)$. Let $W(t)$ be a standard Brownian motion and let $X(t)$ be the Poisson
process with independent increment whose distribution is given by (2) with $\lambda t$ in place of $\lambda$.

1. Theorem. In order that $\bar{N}_{n} \Rightarrow W$ (in the space $D(0, T)$ ), it is necessary and sufficient that for every $\varepsilon>0$ :

$$
\begin{equation*}
\sum_{k:\left|\varepsilon_{k n} / \sigma_{n}\right|>\varepsilon}\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \lambda_{k n} \rightarrow 0 \quad(n \rightarrow \infty) \tag{7}
\end{equation*}
$$

2. Theorem. In order that $N_{n} \Rightarrow X$ (in $D(0, T)$ ) it is necessary and sufficient that, for every $\varepsilon>0$ :

$$
\begin{equation*}
\sum_{k=1}^{p_{n}} \lambda_{k n} \varepsilon_{k n} \rightarrow a+\lambda n \tag{8}
\end{equation*}
$$

$$
\begin{align*}
& \sum_{k:\left|\varepsilon_{k n}-h\right|>\varepsilon} \lambda_{k n} \varepsilon_{k n}^{2} \rightarrow 0 \tag{9}
\end{align*} \quad(n \rightarrow \infty)
$$

## Proof of Theorem 1

Let us first prove the convergence of one dimensional distributions. Without loss of generality set $t=1$ (otherwise we may take $\lambda_{k n} t$ in place of $\lambda_{k n}$ ). By (6), we have to show that (7) is equivalent to

$$
\begin{equation*}
\sum_{k} f\left(\varepsilon_{k n} / \sigma_{n}\right) \cdot\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \lambda_{k n} \rightarrow f(0) \tag{11}
\end{equation*}
$$

as $n \rightarrow \infty$, for all bounded continuous functions on $R$. So, assume (11) and let, for $\varepsilon$ fixed,

$$
f_{\varepsilon}(u)= \begin{cases}1 & \text { if }|u| \geq \varepsilon \\ 0 & \text { if } u=0 \\ 1 \text { inear in }(-\varepsilon, \varepsilon)\end{cases}
$$

Then $\sum_{k:\left|\varepsilon_{k n} / \sigma_{n}\right|>\varepsilon}\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \lambda_{k n} \leq \sum_{k} f_{\varepsilon}\left(\varepsilon_{k n} / \sigma_{n}\right) \cdot\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \lambda_{k n} \rightarrow 0$,
and (7) is proved.
Assume now (7) and let $f$ be a bounded continuous function. So for every $u,|f(u)| \leq M$ and for every $\eta$ there is an $\varepsilon$ such that $f(0)-\eta<f(u)<f(0)+\eta$ if $u \in(-\varepsilon, \varepsilon)$. From (7) it follows $\left|\sum_{k} f\left(\varepsilon_{k n} / \sigma_{n}\right)\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \lambda_{k n}-f(0)\right| \leq n \cdot \sum_{k:\left|\varepsilon_{k n} / \sigma_{n}\right| \leq \varepsilon}\left(\varepsilon_{k n} / \sigma_{n}\right)^{2} \lambda_{k n}+$ $+2 M \sum_{k:\left|\varepsilon_{k n} / \sigma_{n}\right|>\varepsilon}\left(\varepsilon_{\mathrm{k}} / \sigma_{n}\right)^{2} \lambda_{\mathrm{kn}} \rightarrow \eta$
Since $\eta$ is arbitrary small, it proves (11).
The convergence of finite dimensional distributions follows by independence of increments: firstly we have $\left(\bar{N}_{n}\left(t_{1}\right), \bar{N}_{n}\left(t_{2}\right)-\bar{N}_{n}\left(t_{1}\right)\right) \Rightarrow$ $\left(W\left(t_{1}\right), W\left(t_{2}\right)-W\left(t_{1}\right)\right)$ by one-dimensional result; then it follows $\left(\bar{N}_{n}\left(t_{1}\right), \bar{N}_{n}\left(t_{2}\right)\right)=\left(\bar{N}_{n}\left(t_{1}\right), \bar{N}_{n}\left(t_{1}\right)+\left(\bar{N}_{n}\left(t_{2}\right)-\bar{N}_{n}\left(t_{1}\right)\right)\right) \Rightarrow\left(W\left(t_{1}\right), W\left(t_{2}\right)\right)$, and for higher dimension by induction.

Finally, to show tightness, note that, for $t_{1} \leq t \leq t_{2}$ :
$E\left(\left(\bar{N}_{n}(t)-\bar{N}_{n}\left(t_{1}\right)\right)^{2}\left(\bar{N}_{n}\left(t_{2}\right)-\bar{N}_{n}(t)\right)^{2}\right)=E\left(\bar{N}_{n}(t)-\bar{N}_{n}\left(t_{1}\right)\right)^{2} E\left(\bar{N}_{n}\left(t_{2}\right)-\bar{N}_{n}(t)\right)^{2}$

$$
=\left(t-t_{1}\right)\left(t_{2}-t\right) \leq\left(t_{2}-t_{1}\right)^{2}
$$

so by Theorem 15.6. in Billingsley [4], we have that $\bar{N}_{n}$ is tight and $\bar{N}_{\mathrm{n}} \Rightarrow \mathrm{W}$.

Proof of Theorem 2. Note that (9) and (10) together imply

$$
\begin{equation*}
\sum_{k=1}^{P_{n}} \lambda_{k n} \varepsilon_{k n}^{2} \rightarrow \lambda n^{2} \tag{11}
\end{equation*}
$$

Condition (8) implies $G_{n} \Rightarrow G$; so we have to show $K_{n} \Rightarrow K$, i.e., for every bounded continuous $f$ :

$$
\begin{equation*}
\sum f\left(\varepsilon_{k n}\right) \lambda_{k n} \varepsilon_{k n}^{2} \rightarrow f(h) \lambda n^{2} \tag{12}
\end{equation*}
$$

But (12) follows easily from (9) and (11). Assume now (12). Then (9)
and (10) can be proved choosing appropriate $f$ (see the proof of Theorem 1.) This proves one-dimensional convergence. The rest of the proof is same as in Theorem 1. Note that $h$ is determined by (9), $\lambda$ is determined by (10) and then $a$ is determined by (8).

The origin of problems considered here is the following: In the absence of external stimulation, the membrane potential is in the form of impulses of small magnitude arriving at a large number of synaptic sites according to independent Poisson processes (see [10]), so it can be represented by $N_{n}(t)$. If conditions of Theorem 1 are satisfied, then the diffusion approximation is valid, and the membrane potential can be approximated by the solution $V(t)$ of stochastic differential equation.

$$
\begin{equation*}
d V(t)=(-A V(t)+B) d t+\sigma d W(t) \tag{13}
\end{equation*}
$$

where $W$ is the standard Wiener process. Suppose that we want to generalize (13) to infinite dimensions. Then we replace $A$ by a bounded operator; $N_{n}(t)$ becomes a process with independent increments in an infinitely dimensional space. Due to some technical consideration, we may wish to investigate the problem in the Schwartz space of distributions, or on a general multi-Hilbertian space. The convergence results in the previous chapter, the relation between ID laws and homogeneous processes with independent increments give all the necessary techniques needed to extend the example to infinite dimension in exactly similar manner. This was the main motivation for the previous results.

We now investigate purely mathematical questions related to the independent increment processes in the next chapter.

## 4. PROCESSES WITH INDEPENDENT INCREMENTS

### 4.1. Properties of paths

Let $E$ be a countably Hilbert complete space.
Let $(\Omega, F, P)$ be a measure space, and let $X_{t}, 0 \leq t \leq T$ be a random process with independent increments, defined on $\Omega$ and taking values in $E^{\prime}$.

1. Definition. We say that $X_{t}$ is a Lévy process if
(i) $t \rightarrow\left\langle X_{t}, x\right\rangle$ is continuous in probability for every $x$ and
(ii) $X .(\omega) \in D\left([0, T], E^{\prime}\right)$ for almost every $\omega$.
(iii) $X_{0}(\omega)=0$ for almost every $\omega$.

In what follows we shall consider Lévy processes only. Also, we shall assume that (ii) and (iii) hold for every $\omega$, which is not a loss of generality.

From Definition 1, it follows that for each $x,\left\langle X_{t}, x\right\rangle$ is a real valued Lévy process, continuous in probability (see Ito [13] for the definitions on the real line).

By a jump (or a strong jump) of $X(\omega)$ at time $t$ we mean the difference $X_{t}(\omega)-X_{t^{-}}(\omega)$ if it is not zero.

By a weak jump of $X(\omega)$ at time $t$ at the point $x$ we mean
$<x_{t}(\omega)-x_{t}-(\omega)$ if it is not zero.
Clearly, $X$ has a strong jump at $t$ if and only if $X$ has at least one weak jump at t .
2. Lemma. If $x_{s} \rightarrow x$ in $E^{\prime}(a s s \rightarrow t$ ), then there is some $m$ such that $x$ and $X_{s}$ belong to $E_{m}^{\prime}$ for $s$ close enough to $t$, and $x_{s} \rightarrow x$
in $E_{m}^{\prime}$.
Proof. Suppose not. Then, for $x_{s} \rightarrow x$, there exists a sequence $x_{n_{k}}$ such that $x_{n_{k}} \in E_{n_{k}}^{\prime}$ and $n_{k} \rightarrow \infty$. Then $x_{n_{k}} \rightarrow x$, and $x_{n_{k}} \in E_{n_{k}}^{\prime}$ contradicts Theorem 34 in [8], for sequences.
3. Lemma. If $T<\infty$ and $X_{t}(\omega) \in D\left([0, t], E^{\prime}\right)$, then for every $\omega$ there is a norm $q \in \tau$, such that $X_{t}(\omega)$ belongs to $E_{q}^{\prime}$ for all $t$ in $[0, T]$.

Proof. By the Definition 1, for every $t>0$ we have

$$
\lim _{s \rightarrow t} x_{s}^{(\omega)}=x_{t-}^{(\omega)} .
$$

Then by Lemma 2, there is an $\varepsilon_{t}$ and a norm $q_{t}$ such that for $s \in\left(t-\varepsilon_{t}, t\right)$ we have $x_{s}^{(\omega)} \in E_{q_{t}}^{\prime}$. For $t=0$ we have $x_{0}^{(\omega)}=0$, and, by right continuity $X_{s}^{(\omega)} \rightarrow 0^{t}$ as $s \rightarrow 0, s>0$. So, again by Lemma 1 , there is an $\varepsilon_{0}$ such that $X_{t}^{(\omega)}$ belongs to some space $E q_{0}^{\prime}$ for all $x$ in $\left[0, \varepsilon_{0}\right]$.

By compactness, there is a finite cover $0_{i}, 1 \leq i \leq n$, such that all points in $0_{i}$ belong to the space $E q_{k}^{\prime}$. But then all $X_{t}$ for $t \in[0, T]=0_{1} \cup 0_{2} \cup \ldots \cup 0_{n}$ belong to the space $E q^{\prime}=E q^{\prime} \cup \ldots \cup E q_{n}^{\prime}$.
4. Lemma. Suppose $T<\infty$ and suppose that, for every $x \in E$, the number of jumps of $\left\langle\chi_{t}^{(\omega)}, x\right\rangle, 0 \leq t \leq T$ is finite. Then there is a constant $M$ such that for every $x$, the function $\left\langle x_{t}^{(\omega)}, x\right\rangle$ has at most $M$ jumps in $[0, T]$ ( $\omega$ is fixed).

Proof. Let $f_{T}(x)$ be the number of jumps of $\left\langle x_{t}, x\right\rangle$ up to time $T$. For each $x \in E, f_{T}(x)$ is finite. Let us now prove that $f_{T}$ is lower semicontinuous, i.e., if $x_{n} \rightarrow x$ then
(1)

$$
\frac{1 i_{n \rightarrow \infty}}{} f_{t}\left(x_{n}\right) \geq f_{t}(x)
$$

It clearly suffices to show that, whenever $\left\langle X_{t}, x\right\rangle$ has a jump of size $>n$ then, starting with some $n_{0}$ all functions $\left\langle X_{t}, x_{n}\right\rangle$ have a jump of size > $\eta$.

Suppose now that $\langle X, x\rangle$ has a jump of amount $>\eta$ at time $t$.
By assumption $\lim _{s<t} X_{s}=X_{s-}$.

$$
s \rightarrow t
$$

By Lemma 2, there is a norm $q$ such that

$$
\begin{equation*}
\tilde{q}\left(x_{s}-X_{t-}\right)<\varepsilon \quad \text { if } s \in[t-\varepsilon, t) \text { for some } \varepsilon \text {. } \tag{2}
\end{equation*}
$$

Clearly, it suffices to show (1) for $q(x)<1$ only. So, assume $x_{n} \rightarrow x$ and $q\left(x_{n}\right)<1$ and $q(x)<1$. Then we have from (2):

$$
\begin{equation*}
\left|<x_{s}, x_{n}-x_{t-}, x_{n}>\right|<\varepsilon, s \in[t-\varepsilon, t) . \tag{3}
\end{equation*}
$$

By $\tau$-continuity of $X_{t}$ and $X_{t}$ - we have

$$
\begin{align*}
& \left|\left\langle x_{t}, x_{n}\right\rangle-\left\langle x_{t}, x\right\rangle\right|<\varepsilon \quad \text { for } n \geq n_{0}  \tag{4}\\
& \left|\left\langle x_{t-}, x_{n}\right\rangle-\left\langle x_{t-}, x\right\rangle\right|<\varepsilon \quad \text { for } n \geq n_{0} . \tag{5}
\end{align*}
$$

Then from (3), (4) and (5) it follows:

$$
\mid<x_{t-}, x_{n}>-\left\langle x_{t}, x_{n}>\right|>n-3 \varepsilon,
$$

and the assertion (1) follows by arbitrarity of $\varepsilon$.
By Osgood's theorem [12], page 62, , there exists a $\tau$-norm $p$ such that

$$
\left|f_{t}(x)\right|<M / 2
$$

for all $x$ in $B_{r}(z)=\{x: p(z-x) \leq r\}$, for some $z \in E$ and $r>0$.
The function $f_{T}(x)$ has the following properties:

$$
\begin{align*}
& f_{T}(x+y) \leq f_{T}(x)+f_{T}(y)  \tag{6}\\
& f_{T}(c x)=f_{T}(x) \tag{7}
\end{align*}
$$

which follow immediately from the definition of $f_{T}$. So, we deduce:

$$
\begin{aligned}
f_{T}(x) & =f_{T}\left(\left(p^{-1}(x) \cdot r x+z-z\right) \cdot p(x) \cdot r^{-1}\right) \\
& =f_{T}\left(p^{-1}(x) \cdot r x+z-z\right) \\
& \leq f_{T}\left(p^{-1}(x) \cdot r x+z\right)+f_{T}(-z) \\
& \leq M .
\end{aligned}
$$

5. Lemma. If $T<\infty$, then $X_{t}$ has only countably many jumps. Proof. See Remark 3.3.
4.2. A decomposition of $E^{\prime}$-valued processes.

Let $E$ be a countably Hilbert complete space.
Let $E_{0}^{\prime}=E^{\prime}-\{0\}$. Denote by $\mathbb{C}$ the cylinder (= Bore) $\sigma$-algebra of $E^{\prime}$, and by $B$ the Bore $\sigma-a l g e b r a$ of sets in $T_{0} \times E_{0}^{\prime}$, where $T_{0}$ will denote the interval $(0, T]$, for $T$ finite or infinite.

Let $B^{*}$ be the class of all sets $A$ in $B$ such that for some $a>0$ and some $x_{1}, x_{2}, \ldots, x_{k} \in E^{\prime}$ we have

$$
\begin{equation*}
A \subset(0, a) \times\left\{F:\left|<F, x_{1}>\left|>\frac{1}{a}, \ldots,\left|<F, x_{k}>\right|>\frac{1}{a}\right\} .\right.\right. \tag{1}
\end{equation*}
$$

Let $X_{t}$ be a process with independent increments.
Let, for an $\omega$ fixed, $i(\omega)=\left\{t: X_{t} \neq X_{t-}\right\}$ and $\mathcal{J}(\omega)=\left\{\left(t, \Delta X_{t}\right), t \in\right.$ $I(\omega)\}$, where $\Delta X_{t}=X_{t}-X_{t-}$. Define the set function $N$ by

$$
N(A)=N(A, \omega)=\text { number of points in } A \cap J(\omega)
$$

$N(A)$ is a finite random variable for $A \in B^{*}$. The completed cylinder $\sigma$-algebra generated by $X_{u}-X_{v}, s \leq u, v \leq t$ will be denoted by $B_{s t}(X)$.
2. Lemma. If $A \in B^{*}$ and $A \subset(s, t] \times E_{0}^{\prime}$ then $N(A)$ is $\bar{B}_{s t}(X)$-measurable.

Proof. Let $E\left(s, t, x_{1}, \ldots, x_{k}, a_{1}, \ldots, a_{k}\right)=E(s, t, x, a)=$ $(s, t] \times\left\{F \in E^{\prime}:\left\langle F, x_{1}\right\rangle>a_{1}, \ldots,\left\langle F, x_{k}\right\rangle>a_{k}\right\}$,
where $0 \leq s<t<T, a=\left(a_{1}, \ldots a_{k}\right) \in R^{k}, a_{i}>0$ for $1 \leq i \leq k$, $x=\left(x_{1}, \ldots, x_{k}\right) \in E^{k}$.

Let $Q$ be a countable dense subset of ( $s, t]$, including $t$. We have:

$$
\begin{aligned}
& \{N(E(s, t, x, a)) \geq 1\}=\left\{\text { for some } u \in(s, t],\left\langle X_{u}-X_{u}, x_{i}\right\rangle\right. \\
& \left.>a_{i}, i=1, \ldots, k\right\}
\end{aligned}
$$

Now we have

$$
\begin{aligned}
\{N(E(s, t, x, a)) \geq k+1\}= & U_{r \in Q}^{U}(s, t] \\
& \{N(E(s, t, x, a)) \geq k\} \\
& \cap\{N(E(r, t, x, a)) \geq 1\},
\end{aligned}
$$

so, by induction, $\{N(E(s, t, x, a)) \geq k\} \in \bar{B}_{s t}(x)$, thus $N(E(s, t, x, a))$ is $\bar{B}_{s t}(X)$-measurable.

Let $D$ denote the class of all sets $A$ in $B$ such that $N(A \cap E(s, t, x, a))$
is $\bar{B}_{s t}(x)$-measurable. Since $N(A \cap E(x, t, x, a))$ is a bounded measure in A , it follows that $D$ is a Dynkin class. By measurability of $N(E(s, t, x, a))$ we conclude that the class

$$
\begin{gathered}
M=\left\{(s, \infty) \times\left\{F \in E^{\prime}:<F, x_{j}>\geq b_{j}\right\}, j=1, \ldots, n, n \in N,\right. \\
\left.0 \leq s<\infty,-\infty<b_{j}<\infty\right\}
\end{gathered}
$$

belongs to $D$.
$M$ is a multiplicative class that generates $B$; so, by Dynkin's theorem $N(A \cap E(s, t, x, a))$ is $\bar{B}_{s t}(X)$-measurable for every $A$ in $B$.

Now if we define $E^{-}(s, t, x, a)=(s, t] \times\left\{F \in E^{\prime}:\left\langle F, x_{1}\right\rangle \leq-a_{1}, \ldots\right.$, $\left.\left\langle F, x_{k}\right\rangle \leq-a_{k}\right\}$ we similarly conclude that $N\left(A \cap E^{\prime}(s, t, x, a)\right)$ is $\bar{B}_{s t}(X)$-measurable, for every $A$ in $B$.

Finally, if $A \in B^{*}$ and $A \subset(s, t] \times E_{o}^{\prime}$, then for some $a$ and $x$ we have

$$
A=A \cap E(s, t, x, a) \cup A \cap E^{-}(s, t, x, a),
$$

so

$$
N(A)=N(A \cap E(s, t, x, a))+N\left(A \cap E^{-}(s, t, x, a)\right) \text { is }
$$

$\bar{B}_{s t}(X)$-measurable.
Let $B_{C} \subset B$ denote the algebra of all sets of the form: Borel set in $T_{0} \times$ cylinder set in $E_{0}^{\prime}$.
3. Lemma. (i) For $A \in B^{\star}, N(A)$ is Poisson distributed with a finite parameter.
(ii) For $A \in B_{C}, N(A)$ is either identically equal to $\infty$ (a.s.) or is Poisson distributed with finite parameter.

Proof. Let $A(t)=A \cap(0, t] \times E_{0}, A \in B^{*}$.
Let $N(t)=N(t, \omega)=N(A(t), \omega)$.
Clearly, $N(t, i)$ is a right continuous step function in $t$, increasing, with jumps of amount 1. From Lemma 2 and

$$
N(t)-N(s)=N(A(t)-A(s))=N\left(A \cap(s, t] \times E_{0}^{\prime}\right)
$$

it follows that $N(t)$ is a real Lévy process with independent increments. For every $t$ fixed we have

$$
P\left(N(t)-N\left(t^{-}\right) \neq 0\right) \leq P\left(X_{t}-X_{t}-\neq 0\right)=0,
$$

so $N$ is continuous in probability.

So, by results about real processes, $N(t)$ is a Poisson process.
For sufficiently large $t, A=A(t), N(A)=N(A(t))$ and so $N(A)$ is Poisson distributed.

Let $A \in B_{C}$. Then there is an increasing sequence $A_{n} \in B^{*}$ such that $A_{n} \uparrow A$.

Then $N(A)=\lim _{n} N\left(A_{n}\right) . N\left(A_{n}\right)$ is Poisson random variable with expectdion $\lambda_{n}$, say. $\lambda_{n}$ is increasing sequence. If $\lim _{\rightarrow \infty} \lambda_{n}=\infty$, then for every k,

$$
P(N(A) \leq k) \leq P\left(N\left(A_{n}\right) \leq k\right)=e^{-\lambda} n \cdot \sum_{j=0}^{k} \lambda_{n}^{j} / j!\rightarrow 0,
$$

so $N(A)=\infty$ almost surely.

If $\lim _{n} \lambda_{n}=\lambda<\infty$, then $N(A)$ is Poisson with parameter $\lambda$.
Recall that (Section 1), if $X_{t}$ is an E'-valued Levy process with independent increments, then for every $x,\left\langle X_{t}, x\right\rangle$ is a real Lévy process with independent increments. Let $B_{x}\left(T_{0} \times R_{0}\right)$ denote cylinder sets in $T_{0} \times R_{0}$ that depend on $x$ only. If $A \in B_{x}\left(t_{0} \times R_{0}\right)$ then the number of strong jumps of $X_{t}$ that take place in $A$ is equal to the number of jumps of the real process $\left\langle X_{t}, x\right\rangle$ that take place in $A$. Denoting by $N(x, A)$ the number of jumps of $\left\langle X_{t}, x\right\rangle$ in $A$ we have the following:
4. Lemma. If $A \in B_{x}\left(T_{0} \times R_{0}\right)$, then $N(A)=N(x, A)$.

Let $n(A)=E(N(A, \omega))$. Rewriting the Lévy-Ito decomposition for real processes, and using Lemma 3, we have
5. Theorem. For every $x$ in $E\left\langle X_{t}, x\right\rangle=Z_{t}(x)+Y_{t}(x)$, where $Z_{t}(x)$ is a Gaussian process with independent increments, and $Y_{t}(x)$ is given by

$$
\left.\begin{array}{rl}
Y_{t}(x)= & \lim _{k \rightarrow \infty}\left(\int_{0<s \leq t}<F, x>d N(s, F)\right.
\end{array}+\int_{0<s \leq t}<F, x>\operatorname{dn}(s, F)\right)
$$

6. Remark. The decomposition in Theorem 5 is not a decomposition in $E^{\prime}$, because the process $Y_{t}$, and consequently $Z_{t}$ may not be linear in $x$. In the following section we shall prove the complete decomposition in a separable Hilbert space. The general problem remains open.

### 4.3. A Leýy-Ito decomposition

Let $H$ be a real separable Hilbert space. We assume that $X_{t}$ is a $H$ valued Levy process with independent increments, and $t \in[0, T]$, where $T$ is finite or infinite.

1. Theorem. Let $n>0$. If $X_{t} \in D([0, T], H)$ then $X$ has only finitely many jumps of the norm bigger than $\eta$ on any finite subinterval of $[0, T]$.

Proof. Let us first prove the theorem on $[0, T], T<\infty$. Let $A$ be the class of all points in $[0, T]$ such that there are only finitely many jumps of the norm $>\eta$ in $[0, t)$. Let $s=\sup \{t: t \in A\}$. By continuity of $X$ at 0 , and by $X_{0}=0$, there is a neighborhood $U=(0, \varepsilon)$ such that $\left\|X_{t}\right\|<n / 2$ for $t \in U$. Then for $u, v \in U,\left\|X_{u}-X_{v}\right\| \leq\left\|X_{u}\right\|+\left\|X_{v}\right\| \leq n$, so there is no jump in $U$ of size $>\eta$. Thus $s>0$. Suppose $s<1$. Then again, there is a neighborhood ( $s, s+\varepsilon$ ) in which $X$ does not have jumps exceeding $\eta$ in the norm; so $s=T$. The above proof goes through if $[0, T]$ is replaced by an arbitrary finite interval.
2. Theorem. If $X \in D([O, T], H)$, then $X$ has only countably many jumps.

Proof. Immediate by Theorem 1.
3. Remark. If $E$ is a countably Hilbert space, and $X_{t}$ an E'=valued Levy process, it is proved in Lemma 1.2. that for $\omega$ fixed, $X_{t}(\omega)$ belongs to a Hilbert space. So, Theorem 2 holds for $E^{\prime}$.
4. Lemma. Let $F=\left\{B_{s t}, s<t, s, t \in T\right\}$ be an additive family of sub-$\sigma$-algebras on $H$, i.e.,
(i) $B_{s t} \vee B_{t n}=B_{s n}$ a.s.,
(ii) $\left\{B_{t_{0} t_{1}}, B_{t_{1} t_{2}}, \ldots, B_{t_{n-1} t_{n}}\right\}$ is independent for $t_{0}<t_{1} \ldots<t_{n}$.

Let $X_{t}$ be a Lévy process with values in $H, Y_{t}$ a real valued Lévy process of Poisson type, such that with probability one there are no common jump times for $X_{t}$ and $Y_{t}$. Assume that both $X_{t}$ and $Y_{t}$ are adapted to $F$, i.e., $X_{t}-X_{s}$ as well as $Y_{t}-Y_{s}$ is $B_{s t}$-measurable. Then $X_{t}$ and $Y_{t}$ are independent.
Proof. If there are no common jump points for $X_{t}$ and $Y_{t}$, then for every $x$, there are no common jump points for real processes $\left\langle X_{t}\right.$, $x>$ and $Y$. Then by Ito's fundamental lemma $\left\langle X_{t}, x\right\rangle$ and $Y_{t}$ are independent. By a trivial extension of this argument, all finite dimensional processes $\left(\left\langle x_{t}, x_{1}\right\rangle, \ldots,\left\langle x_{t}, x_{k}\right\rangle\right)$ are independent of $\gamma_{t}$, and this gives the result.
5. Lemma. A Lévy process $X_{t}$ whose sample functions $X_{t}(\omega)$ are continuous a.s. is a Gaussian process.

Proof. By finite dimensional case, all finite-dimensional distributions
of $X_{t}-X_{s}$ are Gaussian. Therefore, for every $s, t, X_{t}-X_{s}$ is a Gaussian process.
6. Lemma. If a Lévy process $X_{t}$ is Gaussian, then its sample functions are continuous almost surely.

Proof. Let $X_{t}$ be a Gaussian Lévy process. Then for every $x,\left\langle X_{t}, x\right\rangle$ is a Gaussian process, so $\left\langle X_{t}, x\right\rangle$ is ass. continuous (null sets depend on $t$ ).

By separability of $H$, there is a countable dense set $D=\left\{x_{1}, x_{2}, \ldots\right\}$. Then there is $\Omega_{1} \subset \Omega$, such that for every $\omega \in \Omega_{1}$ and every $x_{i} \in D$, $\left\langle x_{t}(\omega), x_{i}\right\rangle$ is continuous, and $P\left(\Omega_{1}\right)=1$.

Now fix $\omega \in \Omega_{1}$. We want to show that for every $x \in H,\left\langle x_{t}(\omega), x\right\rangle$ is continuous. Since $X_{t}$ is a Lévy process, there is a neighborhood $U$ of $t$ such that $\left\|X_{s}(\omega)\right\|<M$ if $s \in U$, for some $M$. Then, let $\varepsilon>0$ be given, and let $x_{j} \in D$ be such that $\left\|x-x_{j}\right\|<\varepsilon$. Let $s \in U$ be such that $\left\langle X_{s}-X_{t}, x_{j}\right\rangle<\varepsilon$.

Then

$$
\begin{aligned}
&\left|<x_{s}-x_{t}, x>\right| \leq\left|<x_{s}-x_{t}, x_{j}>\left|+\left|<x_{s}-x_{t}, x-x_{j}>\right| \leq\right.\right. \\
& \leq \varepsilon+2 M \varepsilon,
\end{aligned}
$$

which shows continuity of $\left\langle x_{t}, x\right\rangle$.
Now, since for every $x,\left\langle x_{t}(\omega), x\right\rangle$ is continuous, it follows that $X_{t}(\omega)$ is continuous for all $\omega \in \Omega_{1}$.
7. Notation. $I(\omega)=\left\{t: X_{t}(\omega) \neq X_{t}-(\omega)\right\}$

$$
\begin{aligned}
& J(\omega)=\left\{\left(t, \Delta X_{t}(\omega)\right), \quad t \in I(\omega)\right\} \\
& \Delta X_{t}(\omega)=X_{t}(\omega)-X_{t-}(\omega) \\
& T=T-\{0\} ; H_{0}=H-\{0\} .
\end{aligned}
$$

$B=B\left(T_{0} \times H_{0}\right)$ is the class of all Borel subsets of $T_{0} \times H_{0}$. $B^{*}=B^{*}\left(T_{0} \times H_{0}\right)$ is the class of all sets $A \subset B$ such that

$$
A \subset(0, a) \times\left\{F \in H:\|F\|>\frac{1}{a}\right\}, \text { for some } a>0 .
$$

By Theorem 1, for $A \subset B^{*}, A \cap J(\omega)$ is a finite set. Let $N(A)=N(A, \omega)$ denote the number of points in $A \cap J(\omega)$

The proofs of the following two lemmas are almost identical to those of 2.2. and 2.3. and therefore are omitted.
8. Lemma. If $A \in B^{*}\left(T_{0} \times H_{0}\right)$ and $A \subset(s, t] \times H_{0}$, then $N(A)$ is $\bar{B}_{s t}(X)$ measurable.
9. Lemma. For $A \in B\left(T_{0} \times H_{0}\right), N(A)$ is either Poisson distributed or identically $=\infty$.

If $A \in B^{\star}$, then $N(A)$ is Poisson with a finite parameter.
Now consider for every $x \in H$ and $A \in B^{\star}$ :

$$
\begin{aligned}
S(A, x)=S(A, x, \omega)= & \left.\left(t, \Delta \sum_{t}\right) \in A A_{t}\left\{\left\langle X^{( }\right), x\right\rangle-\left\langle X_{t}-(\omega), x\right\rangle\right\} \\
= & \left.\sum(t, F) \in A \xrightarrow{\langle F, x\rangle} \cap\right\}(\omega)
\end{aligned}
$$

Clearly, for $A \in B^{*}$ we have

$$
\begin{gather*}
S(A, x)=\lim _{n \rightarrow \infty} \sum_{k} \frac{k}{n} N\left(A \cap\left\{F:\langle F, x\rangle \in\left(\frac{k-1}{n}, \frac{k}{n}\right]\right\}\right.  \tag{1}\\
=\int_{(t, F) \in A}^{\int}\langle F, x\rangle d N(t, F)
\end{gather*}
$$

Define also $S(A)=S(A, \omega)=\left(t, \Delta \bar{X}_{t} \in A X_{t}\right.$
This sum is finite, so we have $S(A, x)=\langle S(A), x\rangle$.
10. Lemma. Lemma 8 holds for $S(A, x)$ and $S(A)$ in place of $N(A)$.

Proof. For $S(A, x)$ clear from relation (1). Then for $S(A)$ it follows from

$$
S(A)=\sum_{i=1}^{\infty}\left\langle S(A), \mathbf{e}_{\mathbf{i}}>\mathbf{e}_{\mathbf{i}}\right.
$$

where $\left\{e_{i}\right\}_{i=1}^{\infty}$ is an orthonormal basis in $H$.
Set now $\quad A(t)=A \cap[0, t] \times H_{0}, A \in B^{*}$

$$
H(t, A)=X_{t}-S(A(t)) .
$$

11. Lemma. $H(t, A)$ is a Lévy process independent of the process $N(A(t))$. Proof. Immediate, by Lemmas 4, 8 and 9.

The following three lemmas can be proved the same way as for real processes (see [13], Section 1.6.5.
12. Lemma. Let $A_{1}, \ldots, A_{n} \in B^{*}$ be disjoint. Then the following processes are independent:

$$
N\left(A_{1}(t)\right), \ldots, N\left(A_{n}(t)\right), H\left(t, \bigcup_{1}^{n} A_{i}\right) .
$$

13. Lemma. Lemma 12 holds for $S$ in place of $N$.
14. Lemma. Let $A_{1}, \ldots, A_{n}$ be disjoint sets in $B$. Then $N\left(A_{1}\right), \ldots, N\left(A_{n}\right)$ are independent.

Now for $A \in B^{*}$ set

$$
A_{m k}(x)=\left\{(s, F) \in A \cap\left\{F:\langle F, x\rangle \in\left[\frac{k}{m}, \frac{k+1}{m}\right]\right\}\right.
$$

Then

$$
\begin{equation*}
\langle S(A), x\rangle=\int_{A} \int\langle F, x\rangle d N(s, F)=\lim _{m \rightarrow \infty} \sum_{k} \frac{k}{m} N\left(A_{m k}(x)\right) \tag{2}
\end{equation*}
$$

Let $\operatorname{EN}(A)=n(A)$.
15. Lemma. If $A \in B^{*}$, then

$$
E e^{i\langle S(A), x\rangle}=\exp \int_{A} \int\left(e^{i\langle F, x\rangle}-1\right) d n(s, F)
$$

## Proof.

For each $m, k, N\left(A_{m k}(x)\right)$ is Poisson distributed with parameter $n\left(A_{m k}(x)\right)$. So,

$$
E \exp \frac{i k}{m} N\left(A_{m k}(x)\right)=\exp \left(\left(\frac{i k}{m}-1\right) n\left(A_{m k}(x)\right)\right)
$$

By relation (2),

$$
\begin{aligned}
E \exp i<S(A), x\rangle & =\lim _{\rightarrow \infty} E \exp \left(i \sum_{k} \frac{k}{m} N\left(A_{m k}(x)\right)\right) \\
& =\lim _{m \rightarrow \infty} \prod_{k} E \exp \left(\frac{i k}{m} N\left(A_{m k}(x)\right)\right) \\
& =\lim _{m \rightarrow \infty} \prod_{k} \exp \left(\left(e^{\frac{i k}{T M}}-1\right) n\left(A_{m k}(x)\right)\right) \\
& =\lim _{m \rightarrow \infty} \exp \left(\sum\left(e^{\frac{i k}{m}}-1\right) n\left(A_{m k}(x)\right)\right) \\
& \exp \int_{A} \int\left(e^{i<F, x>}-1\right) \operatorname{dn}(s, F) .
\end{aligned}
$$

From the above we obtain
16. Lemma. If $A \in B^{\star}$ is included in $\{(s, F):\|F\|<m\}$ for some $m<\infty$, then

$$
\begin{aligned}
E\langle S(A), x\rangle & =\int_{A} \int\langle F, x\rangle \operatorname{dn}(s, F) \\
\operatorname{Var}\langle S(A), x\rangle & =\int_{A} \int\langle F, x\rangle^{2} d n(s, F)
\end{aligned}
$$

17. Remark. If $(0, t] \times\{F:\|F\|>\varepsilon\} \quad(\varepsilon>0)$ belongs to $B^{\star}\left(T_{0} \times H_{0}\right)$, we have

$$
\begin{aligned}
& \int_{0 \leq \leq \leq t} \int_{\|F\|} d n(s, F)<\infty, \\
& \|F\| \geq \varepsilon
\end{aligned}
$$

because $N$ on sets in $B^{*}$ is Poisson with finite parameter $n$.
18. Lemma. For every $x \in H$,
(i) $\int_{\substack{0<s \leq t}}^{\int F, x>}{ }^{2} d n(s, F)<\infty$
(ii) $\int_{\substack{0<s \leq t \\\|F\| \leq 1}}\|F\|^{2} d n(s, F)<\infty$.

Proof. (i) Let $E_{\varepsilon}=\{(s, F): 0<s \leq t, \varepsilon<\|F\|<1\}, \varepsilon>0$. Then $E_{\varepsilon} \in B^{*}$ and, by Lemma 15,

$$
E e^{i<S\left(E_{\varepsilon}\right), x>}=\exp \int_{\substack{0<s \leq t \\ \varepsilon<\|F\| \leq 1}}\left(e^{i<F, x>}-1\right) \operatorname{dn}(s, F)
$$

By Lemma 13, the processes $S\left(E_{\varepsilon}(s)\right)$ and $X_{s}-S\left(E_{\varepsilon}(s)\right)$ are independent. Thus, $S\left(E_{\varepsilon}\right)$ and $X_{t}-S\left(E_{\varepsilon}\right)$ are independent, and

$$
E e^{i\left\langle X_{t}, x\right\rangle}=E e^{i\left\langle S\left(E_{\varepsilon}\right), x\right\rangle} \cdot E e^{i\left\langle X_{t}-S\left(E_{\varepsilon}\right), x\right\rangle}
$$

So, we have

$$
\begin{aligned}
\left|E e^{i<X_{t}, x>}\right| & \leq\left|E e^{\left.i<S\left(E_{\varepsilon}\right), x\right\rangle}\right| \\
& =\exp \int_{\substack{0<s<t \\
\varepsilon<\|F\| \leq 1}}(\cos \langle F, x>-1) \operatorname{dn}(s, F)
\end{aligned}
$$

If $\|x\| \leq 1$ and $\|F\| \leq 1$, then $|<F, x\rangle \mid \leq 1$; in that case we have $\cos \langle F, x\rangle \leq 1-\langle F, x\rangle^{2} / 4$. So, for every $x$ such that $\|x\| \leq 1$ :

$$
\left|E e^{i\left\langle X_{t}, x>\right.}\right| \leq \exp \left\{\underset{\substack{-\int \leq s \leq t}}{\left.\varepsilon<\|F\| \leq 1 \leq\rangle^{2} \operatorname{dn}(s, F)\right\}}\right.
$$

Letting $\varepsilon+0$, taking logarithm and dividing by $\|x\|$, we have that for every $x \in H$ :

$$
\left.\begin{array}{rl}
\|x\|^{-2} \int_{0<s \leq t} \int & <F, x>^{2} d n(s, F) \tag{3}
\end{array}\right) \leq-\log \left|E e^{i<X_{t}, x>/\|x\|}\right|
$$

where the last inequality follows from the fact (one-dimensional) that characteristic function of $\left\langle x_{t}, x\right\rangle$, being an ID, never vanishes. This proves (i).

The proof of (ii) will be given in Theorem 22.
19. Notations and some facts.

Set

$$
\begin{aligned}
& \begin{aligned}
S_{1}(t, x)=S_{1}(t, x, \omega)= & \int_{\substack{0<s \leq t}}<F, x>d N(s, F) \\
& \|F\| \geq 1
\end{aligned} \\
& =S((0, t] \times\{F:\|F\| \geq 1\}, x) \\
& =\sum_{0<j \leq t}<x_{s}-x_{s-}, x>, \\
& S_{k}(t, x)=S_{k}(t, x, \omega)=\int_{0<s \leq t}\langle F, x\rangle d N(s, F), k>1 . \\
& \frac{1}{k} \leq\|F\|<\frac{1}{k-1} \\
& \text { Define } S_{1}(t)=\sum_{\substack{0<s \leq t \\
\left\|\Delta X_{s}\right\| \geq 1}} \Delta X_{t} \text {, and } S_{k}(t) \text { - analogously. }
\end{aligned}
$$

Clearly, for $k \geq 1, S_{k}(t) \in H$ and $\left\langle S_{k}(t), x\right\rangle=S_{k}(t, x)$.
Define also

$$
\begin{aligned}
T_{k}(t, x)= & T_{k}(t, x, \omega)=S_{1}(t, \omega, x)+\sum_{j=2}^{k} S_{j}(t, x, \omega)-E S_{j}(t, x, \omega) \\
= & \int_{0<s \leq t}<F, x>d N(s, F)-\int_{0<s \leq t}<F, x>\operatorname{dn}(s, F) . \\
& \frac{1}{k} \leq\|F\| \leq 1 / k
\end{aligned}
$$

Using previous results, it can be shown that $S_{n}(t), n=1,2, \ldots$, are independent Lévy processes. By Lemma $18(i), E\left(S_{k}(t, x)\right)$ is continuous in $t$; so $T_{k}(t, x)$ is a Lévy process for every $k$ and $x$.

Using Komogorov's inequality (see [13] Section 1.7.2.) we have the following
20. Theorem. For every $x \in H$, we have the following decomposition:
(4)

$$
\begin{aligned}
\left\langle X_{t}(\omega), x\right\rangle= & Z_{t}(x, \omega)+\lim _{k \rightarrow \infty}\left(\int_{0<s \leq t} \int_{\|F\| \geq 1 / k}<F, x>d N(s, F)\right. \\
& \left.+\int_{0<s \leq t}<F, x>\operatorname{dn}(s, F)\right), \\
& \frac{1}{k} \leq\|F\| \leq 1
\end{aligned}
$$

where, for every $x, Z_{t}(x)$ is a continuous Lévy process (thus Gaussian), and independent on $N$.

Now we can prove the following theorem on a decomposition in Hilbert space.
21. Theorem. Let $X_{t}$ be a Hilbert space-valued Lévy process. Let $N(A, \omega)$ be the number of jumps of $X_{t}(\omega)$ in $A$, and let $n(A)=\operatorname{EN}(A)$. Then $N$ and $n$ are measures on $H$ and we have:

$$
\begin{equation*}
x_{t}=z_{t}+y_{t} \tag{5}
\end{equation*}
$$

where $Z_{t}$ is a Gaussian random process in $H, Y_{t}$ is an $H$-valued random process, independent on $Z_{t}$, defined by

Proof. For $\omega$ and $k$ fixed, the term after the "limp" in (4) is linear in $x$. Notice that

$$
\begin{aligned}
& \int_{0<s \leq t}^{\int}<F, x>d N(s, F) \\
& \|F\| \geq 1 / k
\end{aligned}
$$

is continuous in $x$, as a finite sum of continuous terms.
For the other term, we have

$$
\begin{aligned}
\left|\int_{0<s \leq t}<F, x>d n(s, F)\right| \leq\|x\| \cdot & \int_{0<s \leq t}\|F\| d n(s, F) \\
\frac{1}{k} \leq\|F\| \leq 1 & \frac{1}{k} \leq\|F\| \leq 1 \\
& \leq\|x\| \cdot n\left((0, t] \times\left\{F: \frac{1}{k} \leq\|F\| \leq 1\right),\right.
\end{aligned}
$$

so this is also continuous in $x$.
Then, by Banach-Steinhaus theorem, the limit in (4) defines a continuous linear operator in $H$, call it $Y_{t}(\omega)$.

Thus, $Z_{t}=X_{t}-Y_{t}$ is also in $H$, and, since $\left\langle Z_{t}, x\right\rangle$ is continuous in $t$ for every $x$, it foolows that $Z_{t}$ is also continuous. Then, by Lemma $5, Z_{t}$ is Gaussian on $H . Z_{t}$ and $N$ are independent by Lemma 4.

For a Lévy process $X_{t}$ decomposed as in previous theorem, we define

$$
\begin{aligned}
M_{t}(x) & =E\left\langle Z_{t}, x\right\rangle \\
V_{t}(x) & =\operatorname{Var}\left\langle Z_{t}, x\right\rangle
\end{aligned}
$$

where $Z_{t}$ is Gaussian process defined by (5).
For fixed $t, M_{t}$ is linear and continuous in $x ; V_{t}$ is an $I(\|\cdot\|)$ seminorm.

The functions $M_{t}, V_{t}$, and the measure $d n(s, F)$ are called (following Ito), the three components of $X_{t}$.
22. Theorem. Let $X_{t}$ be a Lévy process with three components $M_{t}, V_{t}, n$.

Then the characteristic function of $x_{t_{2}}-x_{t_{1}}, t_{2}>t_{1}$ is given by

$$
\begin{align*}
f_{t_{1}, t_{2}}(x)= & \exp \left\{i<M_{t_{2}}-M_{t_{1}}, x>-\frac{1}{2}\left(V_{t_{2}}(x)-V_{t_{1}}(x)\right)\right.  \tag{6}\\
& \left.+\int_{\substack{t_{1}<s<t_{2} \\
\\
F \neq 0}}\left(e^{i<F, x>}-1-i<F, x>\cdot 1[\|F\| \leq 1]\right) \operatorname{dn}(s, F)\right\}
\end{align*}
$$

and is subject to an infinitely divisible distributions. Also,

$$
\begin{equation*}
\int_{\substack{0<s \leq t \\\|F\| \leq 1}}\|F\|^{2} \operatorname{dn}(s, F)<\infty \tag{7}
\end{equation*}
$$

Proof. It suffices to prove (6) for $t_{2}=t, t_{1}=0$.
By Theorem 20, we have

$$
\left\langle X_{t}, x\right\rangle=\left\langle Z_{t}, x\right\rangle+\lim _{k \rightarrow \infty} T_{k}(t, x),
$$

and these two terms are independent. So,

$$
\begin{aligned}
E e^{i\left\langle X_{t}, x\right\rangle} & =E e^{i\left\langle Z_{t}, x\right\rangle} \cdot E \lim _{k \rightarrow \infty} e^{i T_{k}(t, x)} \\
& =\exp \left(i\left\langle m_{t}, x\right\rangle-\frac{V(t)}{2}\right) \cdot \lim _{n \rightarrow \infty} E e^{i T_{k}(t, x)}
\end{aligned}
$$

By previous lemmas we have

$$
E e^{\left.i<T_{k}(t, x)\right\rangle}=\exp \left\{\int_{\substack{0<s \leq t \\\|F\| \geq 1 / k}}\left(e^{i<F, x\rangle} \cdot 1\{\|F\| \leq 1\}\right) \operatorname{dn}(s, F)\right\}
$$

which proves (6).
To prove that this is an infinite divisible characteristic function, notice that in the finite-dimensional case we have (ii) of Lemma 18 trivially satisfied, and therefore, by results in Chapter 2, (6) is an ID characteristic function. By Mandrekar-Zinn theorem 2.8. in [18], we conclude that (6) is an ID characteristic functional of a probability measure in H. By Theorem 2.3.14, we have (7).

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