

26. Lévy's Functional Analysis in Terms of an Infinite Dimensional Brownian Motion. I

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§ 1. Introduction. In his book [1], Paul Lévy has extensively developed a potential theory in an infinite dimensional space.

T. Hida and H. Nomoto have constructed the projective limit $(\mathring{S}_\infty, \mu)$ of the topological stochastic family $\{(\mathring{S}_n, \mu_n)\}$ consisting of the open subsets \mathring{S}_n of the finite dimensional spheres S_n and the restrictions μ_n to \mathring{S}_n of the uniform probability measures on S_n such that $\mu_n(\mathring{S}_n) = 1$.

By using this theory, we shall prove the relation :

$$L^2(\mathring{S}_\infty, \mu) = \varprojlim L^2(\mathring{S}_n, \mu_n),$$

and give an interpretation to Lévy's potential theory for Dirichlet problems on the unit ball by introducing the Brownian motion (B, E) on an infinite dimensional space E such that $E \supset \mathring{S}_\infty$. We shall also establish the strong Markov property, the uniform continuity of the paths and the skew product formula of the Brownian motion.

§ 2. Projectively consistent construction of multiple Wiener integrals. First we reformulate T. Hida and H. Nomoto's results [2] in a slightly different manner from theirs. Let S_n be the sphere with center zero and radius $\sqrt{n+1}$ in the $(n+1)$ -dimensional Euclidean space E_{n+1} , and \mathring{S}_n be the open subset of S_n consisting of the points (x_1, \dots, x_{n+1}) :

$$\begin{cases} x_1 = \sqrt{n+1} \prod_{i=1}^n \sin \theta_i, \\ x_k = \sqrt{n+1} \cos \theta_{k-1} \prod_{i=k}^n \sin \theta_i, & (k=2, \dots, n), \\ x_{n+1} = \sqrt{n+1} \cos \theta_n, \end{cases}$$

with the restriction that $(\theta_1, \dots, \theta_n) \in \Pi^n$, where $\Pi^n = \{(\theta_1, \dots, \theta_n); 0 < \theta_1 < 2\pi, 0 < \theta_i < \pi, i=2, \dots, n\}$. We denote by $\pi_{n,m}$ ($n > m$) the projection of \mathring{S}_n to \mathring{S}_m such that the following is commutative :

$$\begin{array}{ccc} \Pi^n \ni (\theta_1, \dots, \theta_n) & \longrightarrow & (\theta_1, \dots, \theta_m) \in \Pi^m \\ \downarrow & & \downarrow \\ \mathring{S}_n & \xrightarrow{\pi_{n,m}} & \mathring{S}_m. \end{array}$$

Set

$$\mathring{S}_\infty = \left\{ x = (x_1, \dots, x_n, \dots); \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=1}^n x_k^2 = 1, x_1 \neq 0 \text{ or } x_2 < 0 \right\}$$

and define a sequence of projections $\{p_n; n \geq 1\}$ by

$$p_n x = x_n \quad \text{for } x = (x_1, \dots, x_n, \dots) \in \dot{S}_\infty.$$

We denote by \dot{S}_∞ the σ -algebra generated by cylinder sets of \dot{S}_∞ and by μ the standard Gaussian white noise on \dot{S}_∞ . Then $\{p_n; n \geq 1\}$ on (\dot{S}_∞, μ) can be viewed as mutually independent random variables which are subject to the standard normal distribution $N(0, 1)$.

Further we define the projection $\pi_n; \dot{S}_\infty \rightarrow \dot{S}_n$ as follows:

$$\pi_n x = \frac{1}{\|x\|_{n+1}}(x_1, \dots, x_{n+1}),$$

where $\|x\|_{n+1}^2 = \frac{1}{n+1} \sum_{k=1}^{n+1} x_k^2$. Then we have the following

Proposition 2.1 (T. Hida and H. Nomoto [2]). *Let \dot{S}_n ($n \geq 1$) be the topological σ -algebra on \dot{S}_n , and μ_n be the restriction to \dot{S}_n of the uniform probability measure on S_n . Then*

- 1) $\bigcup_{n=1}^{\infty} \pi_n^{-1}(\dot{S}_n)$ generates the σ -algebra \dot{S}_∞ ,
- 2) $\mu_n(A) = \mu(\pi_n^{-1}(A))$ for $A \in \dot{S}_n$,
- 3) $\mu_n(\pi_{n,m}^{-1}(A)) = \mu_m(A)$, ($n > m$) for $A \in \dot{S}_m$.

Let \mathcal{H}_n ($1 \leq n < \infty$) be the sequence of the complex Hilbert spaces $L^2(\dot{S}_n, \dot{S}_n, \mu_n)$. We shall define a projection $\rho_{n,m}$ ($n > m$) of \mathcal{H}_n to \mathcal{H}_m using the branching rule of the representation theory of the rotation group $SO(n+1)$ (see [5, pp. 449–451]).

To begin with, for integers $j \geq 2, m \geq k \geq 0$ we put

$$D_{j,k,m}(\theta) = A_{j,k,m} C_{m-k}^{(j-1)/2+k}(\cos \theta)(\sin \theta)^k,$$

where $C_{m-k}^{(j-1)/2+k}$ denotes the Gegenbauer polynomial and the positive constant $A_{j,k,m}$ is determined so as to have

$$\int_0^\pi D_{j,k,m}^2(\theta)(\sin \theta)^{j-1} d\theta = \int_0^\pi (\sin \theta)^{j-1} d\theta.$$

A base of homogeneous harmonic polynomials on \dot{S}_n can be taken to be the family

$$\mathcal{E}_{K_n, \pm}^n(\theta_1, \dots, \theta_n) = e^{\pm i k_1 \theta_1} \prod_{j=2}^n D_{j, k_{j-1}, k_j}(\theta_j),$$

where K_n stands for the sequence of integers (k_1, \dots, k_n) ,

$$0 \leq k_1 \leq k_2 < \dots \leq k_n.$$

Since the system $\bigcup_{K_n} \{\mathcal{E}_{K_n, +}^n, \mathcal{E}_{K_n, -}^n\}$ constitutes a C.O.N.S. in \mathcal{H}_n , we can determine the orthogonal projection $\rho_{n,m}$ ($n > m$) of \mathcal{H}_n to \mathcal{H}_m in terms of $\mathcal{E}_{K_n, \pm}^n$'s:

$$\rho_{n,m} \mathcal{E}_{K_n, \pm}^n = \begin{cases} \mathcal{E}_{K_m, \pm}^m & \text{if } k_m = k_{m+1} = \dots = k_n \\ 0 & \text{otherwise,} \end{cases}$$

where K_m denotes the subsequence (k_1, \dots, k_m) of the given sequence K_n . Thus the projective system $\{\mathcal{H}_n, \rho_{n,m}\}$ has been defined.

Now for an infinite sequence K of integers (k_1, \dots, k_n, \dots) satisfying $0 \leq k_1 \leq \dots \leq k_{n-1} \leq k_n = k_{n+1} = k_{n+2} = \dots$, with an n , we define the

functions $\mathcal{E}_{K,+}, \mathcal{E}_{K,-}$ on \mathring{S}_∞ as follows :

$$\mathcal{E}_{K,\pm}(x) = e^{\pm ik_1\theta_1} \prod_{j=2}^n D_{j,k_{j-1},k_j}(\theta_j) \left(\frac{\Gamma((n+1)/2)}{\Gamma(|K|+(n+1)/2)} \left(\frac{n+1}{2} \right)^{|K|} \right)^{1/2} \|x\|_{n+1}^{|K|}$$

where $|K|$ and $(\theta_1, \dots, \theta_n)$ denote k_n and the Euler angles of $\pi_n x$ for the n , respectively. Since the family $\bigcup_K \{\mathcal{E}_{K,+}, \mathcal{E}_{K,-}\}$ constitutes a C.O.N.S. in the complex Hilbert space $\mathcal{H} = L^2(\mathring{S}_\infty, \mu)$, we can define the orthogonal projection ρ_n of \mathcal{H} to \mathcal{H}_n as follows :

$$\rho_n \mathcal{E}_{K,\pm} = \begin{cases} \mathcal{E}_{K_n,\pm}^n & \text{if } k_n = k_{n+1} = k_{n+2} = \dots \\ 0 & \text{otherwise,} \end{cases}$$

where K_n denotes the subsequence (k_1, \dots, k_n) of the given infinite sequence K . Now, dualizing Proposition 2.1, we have the following

Theorem 2.2. 1) For $f \in \mathcal{H}$,

$$\lim_{n \rightarrow \infty} \int_{S_\infty} |f(x) - (\rho_n f)(\pi_n x)|^2 \mu(dx) = 0.$$

2) Let $\{f_n \in \mathcal{H}_n; n \geq 1\}$ be a projectively consistent sequence, that is, $\rho_{n,m} f_n = f_m$ ($n > m$). Then there exists a function $f \in \mathcal{H}$ such that

$$\rho_n f = f_n \quad (n \geq 1),$$

if and only if $\sup_n \|f_n\|_n < \infty$, where $\|\cdot\|_n$ denotes the norm of \mathcal{H}_n .

§ 3. Infinite dimensional space E and Brownian motion on E .

Set

$$E = \left\{ x = (x_1, \dots, x_n, \dots) \in R^\infty; \sup_n \frac{1}{n} \sum_{k=1}^n x_k^2 < \infty \right\},$$

and introduce semi-norms $\{\|\cdot\|_n; 1 \leq n < \infty\}$:

$$\|x\|_n^2 = \frac{1}{n} \sum_{k=1}^n x_k^2, \quad \|x\|_\infty = \overline{\lim}_{n \rightarrow \infty} \|x\|_n.$$

Let O_1 and O_2 be the topologies induced by the semi-norms $\{\|\cdot\|_n; 1 \leq n < \infty\}$ and the semi-norms $\{\|\cdot\|_n; 1 \leq n < \infty\}$ respectively. The σ -algebra generated by cylinder sets of E will be denoted by \mathcal{E} : $\mathcal{E} = \sigma(p_n; n \geq 1)$, where $p_n x = x_n$ for $x = (x_1, \dots, x_n, \dots) \in E$.

Remark. 1) \mathcal{E} is generated by O_2 -open sets, not by O_1 -open sets.

2) The topological space (E, O_1) is non-separable.

3) For any positive number α and x in E ,

$$\sum_{n=1}^{\infty} \frac{x_n^2}{n (\log n)^{1+\alpha}} < \infty.$$

In the sequel we shall use only the O_1 -topology of E without explicit mentions. Now let $\{w_n(t); n \geq 1\}$ be a family of mutually independent 1-dimensional Wiener processes ($w_n(0) = 0$) on a complete probability space $(\hat{\Omega}, \hat{\mathcal{B}}, \hat{P})$. Then we can prove

Theorem 3.1. 1) For any $x = (x_1, \dots, x_n, \dots) \in E$,

$$\overline{\lim}_{n \rightarrow \infty} \frac{1}{n} \sum_{k=1}^n (x_k + w_k(t))^2 = \|x\|_\infty^2 + t, \quad \text{for any } t \geq 0 \text{ a.s.}$$

$$2) \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=1}^n (w_k(t) - w_k(s))^2 = |t - s| \quad \text{for any } t, s \geq 0 \text{ a.s.}$$

Thus we see that the process $W(t)$

$$W(t) = (w_1(t), \dots, w_n(t), \dots)$$

is living in E , which will be called the E -valued Wiener process. We denote by Ω the totality of continuous sample paths ω on E and by B_t the mapping:

$$B_t(\omega) = \omega(t) \quad \text{for } \omega \in \Omega.$$

Now Theorem 3.1 gives us probability measures P^x on Ω , $x \in E$, such that:

$$P^x(B(t_k) \in A_k, k=1, \dots, n) = \hat{P}(W(t_k) \in A_k - x, k=1, \dots, n)$$

for $0 \leq t_1 < \dots < t_n < \infty$ and $A_1, \dots, A_n \in \mathcal{E}$. Then we have a strong Markov process (Ω, P^x, B_t) with state space (E, \mathcal{E}) . Theorem 3.1 shows

$$(3.1) \quad \|B(t, \omega)\|_\infty^2 = \|B(0, \omega)\|_\infty^2 + t, \quad \text{for any } t \geq 0 \text{ a.s. } P^x, x \in E.$$

Hence denoting by τ the first exit time from the unit ball $D_\infty = \{x \in E; \|x\|_\infty < 1\}$, we have

$$\tau(\omega) = 1 - \|x\|_\infty^2 \quad \text{a.s. } P^x, (x \in D_\infty).$$

We shall call the process $B = (\Omega, P^x, B_t)$ with state space (E, \mathcal{E}) the *Brownian motion on E*.

§ 4. Spherical Brownian motion. In this section we shall see that as in finite dimensional cases, the Brownian motion B on E is factored as the skew product of its radial part and an independent spherical Brownian motion (see [3]). We denote by S_∞ the unit sphere $\{x \in E; \|x\|_\infty = 1\}$, and \tilde{P}^ξ , ($\xi \in S_\infty$) the restriction of the probability measure P^ξ to the set $\tilde{\Omega} = \{\omega \in \Omega; \|\omega(0)\|_\infty = 1\}$. In view of (3.1), by putting

$$\xi_t(\omega) = e^{-t/2} B(e^t - 1, \omega),$$

we have a strong Markov process $(\tilde{\Omega}, \tilde{P}^\xi, \xi_t)$ with state space S_∞ . We also have for any $\xi \in S_\infty$

$$\|\xi(t, \omega) - \xi(s, \omega)\|_\infty^2 = 2(1 - e^{-|t-s|/2}) \quad \text{for any } t, s \geq 0 \text{ a.s. } \tilde{P}^\xi.$$

We shall call this process the *spherical Brownian motion*. Let \dot{P}^r ($r \in (0, \infty)$) be the probability measure on the set $\dot{\Omega} = (0, \infty)$ such that

$$\dot{P}^r(A) = \begin{cases} 1 & \text{if } r \in A \\ 0 & \text{if } r \notin A. \end{cases}$$

We define the mapping r_t as follows: $r_t(\dot{\omega}) = \sqrt{\dot{\omega}^2 + t}$ for $\dot{\omega} \in \dot{\Omega}$. Then we have a deterministic Markov process $(\dot{\Omega}, \dot{P}^r, r_t)$ with state space $(0, \infty)$, which will be called the *radial process*. Associate the *random clock* τ_t :

$$\tau_t(\dot{\omega}) = \int_0^t \frac{ds}{r_s^2(\dot{\omega})} = \log \frac{\dot{\omega}^2 + t}{\dot{\omega}^2}$$

with the radial process. By putting $\hat{E} = \{x \in E; \|x\|_\infty > 0\}$, we have

Theorem 4.1. *The skew product process $(\hat{\Omega}, \hat{P}^x, \hat{B}_t)$ with state space \hat{E} :*

$$\hat{B}_t((\dot{\omega}, \omega)) = r_t(\dot{\omega})\xi(\tau_t(\dot{\omega}), \omega),$$

$$P^x = P^r \times P^\xi \quad \text{for } x = r\xi, r > 0, \xi \in S_\infty, x \in E,$$

is equivalent to the part of the Brownian motion B on the open set \hat{E} .

It seems that this theorem could be identified with the Lévy's formula ([1, p. 305, (5)]) which is expressed in terms of the generators of the Brownian motion, the radial process and the spherical Brownian motion.

Lastly we remark that the spherical Brownian motion $\xi(t)$ on S_∞ is an infinite dimensional Ornstein-Uhlenbeck process. The theorem gives a new approach to investigation of infinite dimensional Brownian motions and Ornstein-Uhlenbeck processes (cf. [4]). More details will be discussed in a forthcoming note.

References

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