Sub-diffusion processes in Hilbert space and their associated stochastic differential equations and Fokker-Planck-Kolmogorov equations

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Abstract. This paper focuses on the time-changed Q-Wiener process, a Hilbert space-valued sub-diffusion. It is a martingale with respect to an appropriate filtration, hence a stochastic integral with respect to it is definable. For the resulting integral, two change of variables formulas are derived. Via a duality theorem for integrals, existence and uniqueness theorems for stochastic differential equations (SDEs) driven by the time-changed Q-Wiener process are discussed. Associated fractional Fokker-Planck-Kolmogorov equations are derived using either a time-changed Itô formula or duality. Connections are established between three integrals driven by time-changed versions of the Q-Wiener process, cylindrical Wiener process, and martingale measure.
Keywords. inverse stable subordinator, time-changed $Q$-Wiener processes, Hilbert space-valued sub-diffusion processes, time-changed stochastic differential equations (SDEs) in Hilbert space, fractional Fokker-Planck-Kolmogorov (FPK) equations.

1. Introduction

The $Q$-Wiener process is a stochastic process with values in a separable Hilbert space which is usually infinite dimensional. Analogous to a classic Brownian motion in finite dimensional space, the $Q$-Wiener process is a Gaussian diffusion process which has independent and stationary increments and a covariance operator $Q [1, 2, 3]$. Stochastic integrals in Hilbert space with respect to the $Q$-Wiener process are constructed as an infinite sum of real-valued stochastic integrals with respect to infinitely many independent standard Brownian motions $[4, 5, 6, 7]$. A Hilbert space of functions is a natural setting for the semigroup approach to deal with stochastic partial differential equations (SPDEs) since the solutions to these equations are typically elements of such a space $[8, 9, 10]$. For motivation and details on the semigroup approach, see $[11, 12]$. From the monographs $[1, 4, 5]$, stochastic differential equations (SDEs) driven by the $Q$-Wiener process in Hilbert space are expressed as

\[
\begin{cases}
  dX(t) = [AX(t) + F(t, X(t))]dt + C dW_t \\
  X(0) = x_0 \in H,
\end{cases}
\]  

(1.1)
where the operator, $A$, is typically an elliptic differential operator and $W_t$ is the $Q$-Wiener process. This type of SDE usually corresponds to stochastic partial differential equations (SPDEs) and the solution, $X(t)$, of the SDE (1.1) is also a diffusion process in Hilbert space. Semigroups and distributions at a fixed time associated with the diffusion processes in Hilbert spaces are studied in [13] and [14]. Furthermore, deterministic Fokker-Planck-Kolmogorov (FPK) equations corresponding to the diffusion process, $X(t)$, in Hilbert space together with the existence and uniqueness of their solutions have been investigated in [15] [16] [17]. Specifically, Bogachev and others have started the study of FPK equations on Hilbert space in [15] [17] and [16] continues this study by proving existence and uniqueness results for irregular, even non-continuous drift coefficients.

Sub-diffusion processes in a finite dimensional space are investigated by Meerschaert and Scheffler [18] [19] and Magdziarz [20]. In particular, the sub-diffusion processes arising as the scaling limit of continuous-time random walks are considered in [18] [19]. The sample path properties of this type of process are investigated by the martingale approach in [20]. It is known that a Brownian motion with an embedded time-change which is the first hitting time process of a stable subordinator of index between 0 and 1 is a sub-diffusion process [21]. Since the stable subordinator process is an increasing Lévy process with jumps, there are intervals for which the inverse process is constant. Therefore, a time-changed Brownian motion process often models phenomena for which there are intervals where the process
does not change \cite{22,23,24}. The transition probabilities, i.e., the densities of the time-changed Brownian motion, satisfy a time-fractional FPK equation appearing in \cite{25}. Furthermore, the stochastic calculus for a time-changed semimartingale and the associated SDEs driven by the time-changed semimartingale are investigated in \cite{26}. See references \cite{27,28,29} as well as the forthcoming book in \cite{30} for further discussion and generalizations of the topics mentioned earlier in this paragraph.

Corresponding sub-diffusion processes in Hilbert space, specifically, the time-changed $Q$-Wiener processes have not been investigated. Also the SDEs driven by a time-changed $Q$-Wiener process and their associated fractional FPK equations have not yet been discussed. In this paper, we focus on the time-changed $Q$-Wiener process and its associated stochastic calculus. Also, SDEs driven by the time-changed $Q$-Wiener process and fractional FPK equations of the solutions to the time-changed SDEs are investigated. Specifically, section 2 introduces the concept of a time-changed $Q$-Wiener process which is a sub-diffusion process. Similar to the time-changed Brownian motion, the time-changed $Q$-Wiener process is proven to be a square-integrable martingale in Hilbert space with respect to an appropriate filtration. Furthermore, the increasing process and the quadratic variation process of the time-changed $Q$-Wiener process are explicitly derived.

Section 3 develops the stochastic integral with respect to the time-changed $Q$-Wiener process. First and second change of variable formulas for the time-changed stochastic integral are provided.
Also the time-changed Itô formula for an Itô process driven by the time-changed $Q$-Wiener process is developed. Moreover, SDEs driven by the time-changed $Q$-Wiener process are introduced. Based on results from Grecksch and Tudor [31] for general martingales in Hilbert spaces, the existence and uniqueness of a mild solution to a type of time-changed SDE is provided. From the duality developed between time-changed SDEs and their corresponding non-time-changed SDEs, the existence and uniqueness of strong solutions to the associated time-changed SDEs are established.

Section 4 discusses a connection between three stochastic integrals, the one driven by a time-changed $Q$-Wiener process in Hilbert space, the one driven by a time-changed cylindrical Wiener process in Hilbert space and the one driven by the martingale measure, known as Walsh’s integral [32] [7]. This connection allows one to potentially interpret mild solutions to SDEs driven by time-changed $Q$-Wiener processes in Hilbert spaces as random field solutions from the Walsh approach to SPDEs by applying the results of Dalang and Quer-Sardanyons [33].

Section 5 derives time-fractional FPK equations associated with the solutions to SDEs driven by time-changed $Q$-Wiener processes in two ways. The first method uses the time-changed Itô formula developed in Section 3. The same time-fractional FPK equations are next derived by applying the duality theorem also developed in Section 3 and the FPK equations associated with the solution of SDEs driven by the classic $Q$-Wiener process.
2. The time-changed $Q$-Wiener process

After recalling the definition of the $Q$-Wiener process in Hilbert space, we introduce the time-changed $Q$-Wiener process in Hilbert space. The related martingale properties of the time-changed $Q$-Wiener process will also be investigated.

**Definition 2.1.** Following [5], let $Q$ be a nonnegative definite, symmetric, trace-class operator on a separable Hilbert space $K$, let $\{f_j\}_{j=1}^{\infty}$ be an orthonormal basis in $K$ diagonalizing $Q$, and let the corresponding eigenvalues be $\{\lambda_j\}_{j=1}^{\infty}$. Let $\{w_j(t)\}_{t \geq 0}, j = 1, 2, \cdots$, be a sequence of independent Brownian motions defined on a filtered probability space $(\Omega, \mathcal{F}, \{\mathcal{F}_t\}_{t \geq 0}, \mathbb{P})$. Then the process

$$W_t := \sum_{j=1}^{\infty} \lambda_j^{1/2} w_j(t) f_j$$

is called a $Q$-Wiener process in $K$.

For more details of properties of $Q$-Wiener processes, see [4] [5].

Before introducing the time-changed $Q$-Wiener processes, we need to introduce the time-change process which is used throughout this paper. The time-change applied in this paper is the first hitting time process of a $\beta$-stable subordinator defined as

$$E_t := E_\beta(t) = \inf\{\tau > 0 : U_\beta(\tau) > t\}, \quad (2.1)$$

where $U_\beta(t)$ is the $\beta$-stable subordinator which has index $\beta \in (0, 1)$ and Laplace transform

$$\mathbb{E}(e^{-u U_\beta(\tau)}) = e^{-\tau u^\beta}. \quad (2.2)$$
Note that $E_t$ is also called the inverse $\beta$-stable subordinator. Let $B(t)$ denote a one-dimensional standard Brownian motion. Consider $Z_\beta(t) := B(E_t)$, a subordinated Brownian motion which has been time-changed by $E_t$. The following result of Magdziarz [20] shows that $Z_\beta(t)$ is a square integrable martingale with respect to the appropriate right-continuous filtration,

$$\bar{\mathcal{F}}_t = \bigcap_{u > t} \{ \sigma[B(s) : 0 \leq s \leq u] \lor \sigma[E_s : s \geq 0] \}, \quad (2.3)$$

where $\bar{\mathcal{F}}_0$ is assumed to be complete.

**Theorem 2.2.** Magdziarz, [20] The time-changed Brownian motion $Z_\beta(t)$ is a mean zero and square integrable martingale with respect to the filtration $\{\bar{\mathcal{F}}_{E_t}\}_{t \geq 0}$. The quadratic variation process of $Z_\beta(t)$ is $\langle Z_\beta(t), Z_\beta(t) \rangle = E_t$.

The time-changed Brownian motion, $Z_\beta(t)$, is a sub-diffusion process. By incorporating the time-change, $E_t$, into the independent Brownian motions in Definition 2.1, we define a Hilbert space-valued time-changed $Q$-Wiener process as follows:

**Definition 2.3.** With $Q$, $\{f_j\}_{j=1}^\infty$, and $\{\lambda_j\}_{j=1}^\infty$ as defined in 2.1 let $\{w_j(t)\}_{t \geq 0}, j = 1, 2, \cdots$, be a sequence of independent Brownian motions defined on $(\Omega, \mathcal{F}, \{\mathcal{F}_t\}_{t \geq 0}, \mathbb{P})$ which are independent of $E_t$. Then the process

$$W_{E_t} := \sum_{j=1}^\infty \lambda_j^{1/2} w_j(E_t)f_j \quad (2.4)$$

is called a time-changed $Q$-Wiener process in $K$. 
This time-changed $Q$-Wiener process can be considered as an infinite-dimensional Hilbert space analog of a one-dimensional time-changed Brownian motion. It is a sub-diffusion process in Hilbert space. Let $\mu_t$ be the Borel probability measure induced by the time-changed $Q$-Wiener process, $W_{E_t}$ on $K$, i.e., $\mathbb{E}(W_{E_t}) = \int_K x \mu_t(dx)$.

Then the time-fractional FPK equation corresponding to the time-changed $Q$-Wiener process in the following theorem can be considered as a special case of Theorem 5.3 in Section 5.

**Theorem 2.4.** Suppose $\mu_t$ is the probability measure induced by the time-changed $Q$-Wiener process, $W_E$, on $K$. Then $\mu_t$ satisfies the following time-fractional PDE

$$D_t^{\beta} \mu_t = D_x^2 \mu_t,$$

where $D_x^2$ denotes the second-order Fréchet derivative in space and $D_t^{\beta}$ is the Caputo time fractional derivative operator defined as

$$D_t^{\beta} f(t) = \frac{1}{\Gamma(1 - \beta)} \int_0^t \frac{f'(\tau)}{(t - \tau)^\beta} d\tau,$$

where $\beta$ is the index associated with the $\beta$-stable subordinator $U_\beta(t)$ and $\Gamma(\beta)$ is the gamma function.

The time-changed $Q$-Wiener process is our main object of study. In order to define integrals with respect to this process, and, ultimately, to consider SDEs driven by this process, it is advantageous to view the time-changed $Q$-Wiener process as a martingale with respect to an appropriate filtration. To prove this, we begin with the definition of a martingale in a Hilbert space.
Definition 2.5. Following [5], let $K$ be a separable Hilbert space endowed with its Borel $\sigma$-field $\mathcal{B}(K)$. Fix $T > 0$ and let $(\Omega, \mathcal{F}, \{\mathcal{F}_t\}_{t \leq T}, \mathbb{P})$ be a filtered probability space and $\{M_t\}_{t \leq T}$ be a $K$-valued process adapted to the filtration $\{\mathcal{F}_t\}_{t \leq T}$. Assume that $M_t$ is integrable, i.e., $\mathbb{E}\|M_t\| < \infty$.

a) If for any $0 \leq s \leq t$, $\mathbb{E}(M_t|\mathcal{F}_s) = M_s$, $\mathbb{P}$-a.s., then $M_t$ is called an $\mathcal{F}_t$-martingale.

b) If $\mathbb{E}\|M_T\|^2 < \infty$, the martingale $M_t$ is called square integrable on $0 \leq t \leq T$.

The condition for $\{M_t\}_{t < T}$ to be a martingale in Hilbert space is equivalent to the following: for all $s < t$ and $h \in K$,

$$
\mathbb{E}(\langle M_t, h \rangle | \mathcal{F}_s) = \langle M_s, h \rangle, \quad \mathbb{P} - a.s.
$$

Let $\mathcal{M}_T^2(K)$ denote the collection of all continuously square integrable martingales in the Hilbert space $K$.

Example 1. The $Q$-Wiener process $W_t := \sum_{j=1}^{\infty} \lambda_j^{1/2} w_j(t)f_j$ introduced in Definition 2.1 is a square integrable martingale with respect to the filtration $\mathcal{F}_t$ generated by independent Brownian motions, i.e., $\sigma(w_j(s) : 0 \leq s \leq t, j = 1, 2, \ldots)$.

Define

$$
\tilde{\mathcal{F}}_t = \bigcap_{u > t} \{ \sigma[w_j(s) : 0 \leq s \leq u, j = 1, 2, \ldots] \vee \sigma[E_s : s \geq 0] \}, \quad (2.5)
$$

where $\tilde{\mathcal{F}}_0$ is assumed to be complete, and $\{w_j(t)\}_{t \geq 0}, j = 1, 2, \ldots$, is a sequence of independent Brownian motions which are independent.
of $E_t$. The following theorem establishes that the time-changed $Q$-
Wiener process is also a square-integrable martingale, i.e., $W_{E_t} \in \mathcal{M}^2_T(K)$ with respect to the filtration $\tilde{F}_{E_t}$.

**Theorem 2.6.** The time-changed $Q$-Wiener process defined by Definition 2.3 is a $K$-valued square integrable martingale with respect to the filtration $\mathcal{G}_t := \tilde{F}_{E_t}$.

**Proof.** From [20], all moments of the time change, $E_t$, are finite, i.e.,

$$\mathbb{E}(E_t^n) = \frac{t^{n\beta}n!}{\Gamma(n\beta + 1)},$$

for $n = 1, 2, \ldots$. Let $f_{E_t}$ be the density function for $E_t$. Then, the second moment for the time-changed Brownian motion $w_j(E_t)$ is given by

$$\mathbb{E}(w_j^2(E_t)) = \int_0^{\infty} \mathbb{E}(w_j^2(\tau)) f_{E_t}(\tau) d\tau = \int_0^{\infty} \tau f_{E_t}(\tau) d\tau = \frac{t^{\beta}}{\Gamma(\beta + 1)}.$$

Thus,

$$\mathbb{E}\|W_{E_t}\|_K^2 = \mathbb{E}\left\langle \sum_{j=1}^{\infty} \lambda_j^{1/2} w_j(E_t)f_j, \sum_{i=1}^{\infty} \lambda_i^{1/2} w_i(E_t)f_i \right\rangle_K = \mathbb{E} \sum_{j=1}^{\infty} \lambda_j w_j^2(E_t)$$

$$= \sum_{j=1}^{\infty} \lambda_j \mathbb{E}(w_j^2(E_t)) = \frac{t^{\beta}}{\Gamma(\beta + 1)} \sum_{j=1}^{\infty} \lambda_j < \infty.$$

The sum is finite since $Q$ is a trace-class operator. Also, since the $Q$-
Wiener process, $W_t$, is a square integrable martingale in the Hilbert
space $K$, it follows that for any $h \in K$, the process $X_t$ defined by

$$X_t := \langle W_t, h \rangle_K$$
is a real-valued square integrable martingale with respect to the filtration \( \tilde{F}_t \). This means that in order to prove the time-changed \( \tilde{Q} \)-Wiener process, \( W_{E_t} \), is a square integrable martingale, it suffices to verify that the time-changed real-valued process, \( X_{E_t} \), defined by

\[
X_{E_t} := \langle W_{E_t}, h \rangle_K,
\]

is a square integrable martingale with respect to the filtration \( \tilde{G}_t = \tilde{F}_{E_t} \). Define the sequence of \( \{\tilde{F}_\tau\}\)-stopping times, \( T_n \), by

\[
T_n = \inf\{\tau > 0 : |X(\tau)| \geq n\}.
\]

It is known that the stopped process \( X(T_n \wedge \tau) \) is a bounded martingale with respect to \( \tilde{F}_\tau \). Thus, by Doob’s Optional Sampling Theorem, for \( s < t \),

\[
\mathbb{E}(X(T_n \wedge E_t) \mid \tilde{G}_s) = X(T_n \wedge E_s) \quad (2.6)
\]

The right hand side of (2.6) converges to \( X(E_s) \) as \( n \to \infty \). For the left hand side,

\[
|X(T_n \wedge E_t)| \leq \sup_{0 \leq s \leq t} |X(E_s)|.
\]
Then, applying Doob’s Maximal Inequality yields

\[
\mathbb{E}(\sup_{0 \leq s \leq t} |X(E_s)|^2) = \mathbb{E}(\sup_{0 \leq s \leq E_t} |X(s)|^2) = \int_0^\infty \mathbb{E}(\sup_{0 \leq s \leq \tau} |X(s)|^2 \mid E_t = \tau)f_{E_t}(\tau) \, d\tau
\]

\[
= \int_0^\infty \mathbb{E}(\sup_{0 \leq s \leq \tau} |X(s)|^2)f_{E_t}(\tau) \, d\tau \leq 4 \int_0^\infty \mathbb{E}(|X(\tau)|^2)f_{E_t}(\tau) \, d\tau
\]

\[
= 4 \int_0^\infty \mathbb{E}(|\langle W\tau, h \rangle_K|^2)f_{E_t}(\tau) \, d\tau \leq 4 \int_0^\infty \mathbb{E}(\|W\tau\|_K^2 \|h\|_K^2)f_{E_t}(\tau) \, d\tau
\]

\[
= 4\|h\|_K^2 \int_0^\infty \mathbb{E}\|W\tau\|_K^2 \|W_{E_t}\|_K^2 < \infty. \]

Also by Holder’s inequality,

\[
\mathbb{E}(\sup_{0 \leq s \leq t} |X_{E_s}|) \leq (\mathbb{E}(\sup_{0 \leq s \leq t} |X_{E_s}|^2)^{1/2}.
\]

Thus,

\[
\mathbb{E}(\sup_{0 \leq s \leq t} |X_{E_s}|) < \infty.
\]

By the dominated convergence theorem,

\[
\mathbb{E}(X(T_n \wedge E_t) \mid \mathcal{G}_s) \longrightarrow \mathbb{E}(X(E_t) \mid \mathcal{G}_s), \quad \text{as } n \to \infty.
\]

Therefore, from (2.6),

\[
\mathbb{E}(X(E_t) \mid \mathcal{G}_s) = X(E_s),
\]

which implies that \(X(E_t) = \langle W_{E_t}, h \rangle_K\) is a martingale with respect to the filtration \(\tilde{\mathcal{F}}_{E_t}\). Therefore, \(W_{E_t}\) is a square integrable martingale in the Hilbert space \(K\). \(\Box\)
The following definition and lemma concern the increasing processes and the quadratic variation processes of square integrable martingales in a Hilbert space. These will be applied in Proposition 2.9 to the time-changed $Q$-Wiener process $W_{E_t}$.

**Definition 2.7.** Following [5], let $M_t \in \mathcal{M}_T^2(K)$. Denote by $\langle M \rangle_t$ the unique adapted continuous increasing process starting from 0 such that $\|M_t\|^2_K - \langle M \rangle_t$ is a continuous martingale. The quadratic variation process $\langle \langle M \rangle \rangle_t$ of $M_t$ is an adapted continuous process starting from 0, with values in the space of nonnegative definite trace-class operators on $K$, such that for all $h, g \in K$,

$$
\langle M_t, h \rangle_K \langle M_t, g \rangle_K - \langle \langle \langle M \rangle \rangle_t(h), g \rangle_K
$$

is a martingale.

**Lemma 2.8.** Following [5], the quadratic variation process of a martingale $M_t \in \mathcal{M}_T^2(K)$ exists and is unique. Moreover,

$$
\langle M \rangle_t = tr(\langle \langle M \rangle \rangle_t).
$$

**Proposition 2.9.** The increasing process and quadratic variation process of the time-changed $Q$-Wiener process in Definition 2.3 are respectively

$$
\langle W_E \rangle_t = tr(Q)E_t \text{ and } \langle \langle W_E \rangle \rangle_t = QE_t.
$$

**Proof.** Let $Q$ be a nonnegative definite, symmetric, trace-class operator on a separable Hilbert space $K$ and let $\{f_j\}_{j=1}^\infty$ be an orthonormal basis in $K$ diagonalizing $Q$ with corresponding eigenvalues $\{\lambda_j\}_{j=1}^\infty$. ...
Then,

\[ \| W_{E_t} \|_K^2 = \left\langle \sum_{j=1}^\infty \lambda_j^{1/2} w_j(E_t) f_j, \sum_{i=1}^\infty \lambda_i^{1/2} w_i(E_t) f_i \right\rangle_K = \sum_{j=1}^\infty \lambda_j w_j^2(E_t). \]

On the other hand,

\[ tr(Q)E_t = E_t \sum_{j=1}^\infty \lambda_j. \]

So define the process, \( N_{E_t} \), as

\[ N_{E_t} = \| W_{E_t} \|_K^2 - tr(Q)E_t = \sum_{j=1}^\infty \lambda_j (w_j^2(E_t) - E_t), \]

which can be considered as a time-change of \( N_t \) where

\[ N_t = \| W_t \|_K^2 - tr(Q)t = \sum_{j=1}^\infty \lambda_j (w_j^2(t) - t). \]

From Definition 2.7, \( N_t \) is a real-valued martingale since \( tr(Q)t \) is the unique increasing process of the \( Q \)-Wiener process. By an argument similar to that of Theorem 2.6, the time-changed process, \( N_{E_t} \), is a martingale. Further, \( W_{E_t} \) is a martingale and there is an increasing process, \( \langle W_{E} \rangle_t \), such that \( \| W_{E_t} \|_K^2 - \langle W_{E} \rangle_t \) is a martingale. Finally, since \( \| W_{E_t} \|_K^2 \) is a real-valued submartingale, by the uniqueness of the Doob-Meyer decomposition [34],

\[ \langle W_{E} \rangle_t = tr(Q)E_t. \] (2.7)

Again, by Theorem 2.6 and Lemma 2.8, the quadratic process, \( \langle \langle W_{E} \rangle \rangle_t \), of the time-changed \( Q \)-Wiener process, \( W_{E_t} \), exists and is unique,
and satisfies

\[ tr(\langle\langle W_E \rangle\rangle_t) = \langle W_E \rangle_t. \]  \hspace{1cm} (2.8)

Therefore, from (2.7) and (2.8),

\[ \langle\langle W_E \rangle\rangle_t = QE_t. \]

\[ \square \]

3. SDEs driven by the time-changed $Q$-Wiener process

In this section, we begin by developing the Itô stochastic integral with respect to the time-changed $Q$-Wiener process in Hilbert space. Also a time-changed Itô formula for an Itô process driven by a time-changed $Q$-Wiener process is developed. Finally, the existence and uniqueness of solutions to the time-changed Hilbert space-valued SDEs are investigated.

3.1. Stochastic integral with respect to the time-changed $Q$-Wiener process

In order to construct an Itô stochastic integral with respect to the time-changed $Q$-Wiener process, we briefly recall Itô stochastic integrals with respect to a $Q$-Wiener process without a time change as in [4, 5].

As in Section 2 let $K$ and $H$ be two separable Hilbert spaces, and $Q$ be a symmetric, nonnegative definite trace-class operator on $K$. Let $\{f_j\}_{j=1}^\infty$ be an orthonormal basis (ONB) in $K$ such that $Qf_j = \lambda_j f_j$, where these eigenvalues $\lambda_j > 0$ for all $j = 1, 2, \ldots$. Then the
separable Hilbert space $K_Q = Q^{1/2}K$ with an ONB $\{\lambda_j^{1/2}f_j\}_{j=1}^{\infty}$ is endowed with the following scalar product

$$\langle u, v \rangle_{K_Q} = \sum_{j=1}^{\infty} \frac{1}{\lambda_j} \langle u, f_j \rangle_K \langle v, f_j \rangle_K.$$ 

Let $\mathcal{L}_2(K_Q, H)$ be the space of Hilbert-Schmidt operators from $K_Q$ to $H$. The Hilbert-Schmidt norm of an operator $L \in \mathcal{L}_2(K_Q, H)$ is given by

$$\|L\|_{\mathcal{L}_2(K_Q, H)}^2 = \|LQ^{1/2}\|_{\mathcal{L}_2(K, H)}^2 = tr((LQ^{1/2})(LQ^{1/2})^*).$$

The scalar product between two operators $L, M \in \mathcal{L}_2(K_Q, H)$ is defined by

$$\langle L, M \rangle_{\mathcal{L}_2(K_Q, H)} = tr((LQ^{1/2})(MQ^{1/2})^*).$$

Define $\Lambda_2(K_Q, H)$ as the class of $\mathcal{L}_2(K_Q, H)$-valued processes which are measurable mappings from

$$([0, T] \times \Omega, \mathcal{B}([0, T]) \times \mathcal{F})$$

to

$$\left(\mathcal{L}_2(K_Q, H), \mathcal{B}(\mathcal{L}_2(K_Q, H))\right),$$

adapted to the filtration $\{\mathcal{F}_t\}_{t \leq T}$, and satisfying the condition

$$\mathbb{E} \int_0^T \|\Phi(t)\|_{\mathcal{L}_2(K_Q, H)}^2 dt < \infty.$$
Note that $\Lambda_2(K_Q, H)$ is a Hilbert space if it is equipped with the norm

$$
\|\Phi\|_{\Lambda_2(K_Q, H)} := \left( \mathbb{E} \int_0^T \|\Phi(t)\|_{L_2(K_Q, H)}^2 dt \right)^{1/2}.
$$

The following lemma from [5] can be considered as a definition of a stochastic integral with respect to the $Q$-Wiener process:

**Lemma 3.1.** Mandrekar, [5] Let $W_t$ be a $Q$-Wiener process in a separable Hilbert space $K$, $\Phi \in \Lambda_2(K_Q, H)$, and $\{f_j\}_{j=1}^\infty$ be an ONB in $K$ consisting of eigenvectors of $Q$. Then,

$$
\int_0^t \Phi(s) dW_s = \sum_{j=1}^\infty \int_0^t \left( \Phi(s) \lambda_j^{1/2} f_j \right) d\langle W_s, \lambda_j^{1/2} f_j \rangle_{K_Q}.
$$

In order to incorporate the time-change, $E_t$, into the Itô stochastic integral, the generalized $\tilde{\Lambda}_2(K_Q, H)$ is also considered as the class of $L_2(K_Q, H)$-valued processes which are measurable mappings from

$$
([0,T] \times \Omega, \mathcal{B}([0,T]) \times \bar{\mathcal{F}}_t)
$$

to

$$
(L_2(K_Q, H), \mathcal{B}(L_2(K_Q, H))),
$$
adapted to the filtration $\{\bar{\mathcal{F}}_{E_t}\}_{t \leq T}$, and satisfying the condition

$$
\mathbb{E} \int_0^T \|\Phi(t)\|_{L_2(K_Q, H)}^2 dE_t < \infty.
$$

Similarly, $\tilde{\Lambda}_2(K_Q, H)$ is a separable Hilbert space if it is equipped with the norm

$$
\|\Phi\|_{\tilde{\Lambda}_2(K_Q, H)} := \left( \mathbb{E} \int_0^T \|\Phi(t)\|_{L_2(K_Q, H)}^2 dE_t \right)^{1/2}.
$$
Thus, Itô stochastic integrals with respect to the time-changed $Q$-Wiener process can be introduced.

**Definition 3.2.** Let $W_{E_t}$ be a time-changed $Q$-Wiener process in a separable Hilbert space $K$, $\Phi \in \tilde{\Lambda}_2(K_Q, H)$, and let $\{f_j\}_{j=1}^{\infty}$ be an ONB in $K$ consisting of eigenvectors of $Q$. Then, 

$$
\int_0^t \Phi(s)dW_{E_s} = \sum_{j=1}^{\infty} \int_0^t (\Phi(s)\lambda_j^{1/2}f_j)d\langle W_{E_s}, \lambda_j^{1/2}f_j \rangle_{K_Q}.
$$

Now that the Itô integral with respect to the time-changed $Q$-Wiener process has been established, the next step is to derive the Itô isometry first for elementary processes, and then by extension, for arbitrary processes in $\tilde{\Lambda}_2(K_Q, H)$. Consider the class of $\{G_t\}$-adapted elementary processes of the form

$$
\Phi(t, \omega) = \phi(\omega)1_{\{0\}}(t) + \sum_{j=0}^{n-1} \phi_j(\omega)1_{(t_j, t_{j+1}]}(t),
$$

where $0 \leq t_0 \leq t_1 \leq \cdots \leq t_n = T$ and $\phi, \phi_j, j = 0, 1, \cdots, n - 1$ are respectively $G_0$-measurable and $G_{t_j}$-measurable $L_2(K_Q, H)$-valued random variables such that $\phi(\omega), \phi_j(\omega)$ are linear, bounded operators from $K$ to $H$. Let $\mathcal{E}(\mathcal{L}(K, H))$ denote this class of elementary processes. Proceeding to the Itô isometry for an elementary process $\Phi(t, \omega)$, we need the following useful lemma.

**Lemma 3.3.** Let $\{f_j\}_{j=1}^{\infty}$ be an ONB in $K$ consisting of eigenvectors of $Q$ and $G_s = \tilde{F}_{E_s}$ be the filtration. Then, for $l \neq l'$ and $t > s > 0$,

$$
\mathbb{E}\left(\mathbb{E}\left(\left\langle W_{E_t} - W_{E_s}, f_l \right\rangle_K \left\langle W_{E_t} - W_{E_s}, f_{l'} \right\rangle_K \bigg| G_s\right)\right) = 0.
$$
Proof. From Definition 2.3 in Section 2

\[ W_{E_t} = \sum_{j=1}^{\infty} \lambda_j^{1/2} w_j(E_t)f_j. \]

Using the definition of \( W_{E_t} \),

\[
\mathbb{E}\left( \mathbb{E}\left( \langle W_{E_t} - W_{E_s}, f_t \rangle_K \mathbb{E}\left( W_{E_t} - W_{E_s}, f_{t'} \rangle_K \mid \mathcal{G}_s \right) \right) \right)
\]

\[
= \mathbb{E}\left\{ \mathbb{E}\left( \langle W_{E_t}, f_t \rangle_K \mathbb{E}\left( W_{E_t}, f_{t'} \rangle_K - \langle W_{E_s}, f_t \rangle_K \mathbb{E}\left( W_{E_t}, f_{t'} \rangle_K \mid \mathcal{G}_s \right) \right) \right\}
\]

\[
= \mathbb{E}\left( \mathbb{E}\left( \lambda_t^{1/2} w_t(E_t)\lambda_{t'}^{1/2} w_{t'}(E_t) \mid \mathcal{G}_s \right) \right) - \mathbb{E}\left( \mathbb{E}\left( \lambda_t^{1/2} w_t(E_t)\lambda_{t'}^{1/2} w_{t'}(E_t) \mid \mathcal{G}_s \right) \right)
\]

\[
= \mathbb{E}\left( \mathbb{E}\left( \lambda_t^{1/2} w_t(E_s)\lambda_{t'}^{1/2} w_{t'}(E_t) \mid \mathcal{G}_s \right) \right) + \mathbb{E}\left( \lambda_t^{1/2} w_t(E_s)\lambda_{t'}^{1/2} w_{t'}(E_s) \right)
\]

\[ := I_1 - I_2 - I_3 + I_4. \]

Since \( w_t \) is independent of \( w_{t'} \), conditioning on \( E_t \) to compute the first term yields

\[ I_1 = \mathbb{E}\left( \mathbb{E}\left( \lambda_t^{1/2} w_t(E_t)\lambda_{t'}^{1/2} w_{t'}(E_t) \mid \mathcal{G}_s \right) \right) = \mathbb{E}(\lambda_t^{1/2} w_t(E_t)\lambda_{t'}^{1/2} w_{t'}(E_t)) \]

\[ = \lambda_t^{1/2} \lambda_{t'}^{1/2} \int_0^{\infty} \mathbb{E}(w_t(\tau)w_{t'}(\tau)f_{E_t}(\tau))d\tau = \lambda_t^{1/2} \lambda_{t'}^{1/2} \int_0^{\infty} 0 \cdot f_{E_t}(\tau)d\tau = 0. \]

On the other hand,

\[ I_2 = \mathbb{E}\left( \mathbb{E}\left( \lambda_t^{1/2} w_t(E_s)\lambda_{t'}^{1/2} w_{t'}(E_t) \mid \mathcal{G}_s \right) \right) = \lambda_t^{1/2} \lambda_{t'}^{1/2} \mathbb{E}\left( \mathbb{E}\left( w_t(E_s)w_{t'}(E_t) \mid \mathcal{G}_s \right) \right) \]

\[ = \lambda_t^{1/2} \lambda_{t'}^{1/2} \mathbb{E}\left( w_t(E_s) \left( w_{t'}(E_t) - w_{t'}(E_s) \right) + w_t(E_s)w_{t'}(E_s) \mid \mathcal{G}_s \right) \]

\[ = \lambda_t^{1/2} \lambda_{t'}^{1/2} \mathbb{E}\left( w_t(E_s) \left( w_{t'}(E_t) - w_{t'}(E_s) \right) \mid \mathcal{G}_s \right) + w_t(E_s)w_{t'}(E_s) \]

\[ = 0, \]
since $\mathbb{E}(w_t(E_t) - w_t(E_s) | G_s) = 0$ by the martingale property of $W_{E_t}$ and $\mathbb{E}(w_t(E_s) w_t(E_s)) = 0$ by the same conditioning argument previously used in computing term $I_1$. Similarly, the third term, $I_3$, and the fourth term, $I_4$, are also equal to 0. □

**Theorem 3.4.** Let $\Phi \in \mathcal{E}(\mathcal{L}(K, H))$ be a bounded elementary process. Then, for $t \in [0, T]$,

$$
\mathbb{E} \left\| \int_0^t \Phi(s)dW_{E_s} \right\|_H^2 = \mathbb{E} \int_0^t \| \Phi(s) \|_{L_2(K_Q, H)}^2 dE_s < \infty.
$$

**Proof.** First, without loss of generality, assume that $t = T$. Then, for the bounded elementary process, $\Phi$, defined in (3.1),

$$
\mathbb{E} \left\| \int_0^T \Phi(s)dW_{E_s} \right\|_H^2 = \mathbb{E} \left\| \sum_{j=0}^{n-1} \phi_j(W_{E_{t_{j+1}}} - W_{E_{t_j}}) \right\|_H^2
$$

$$
= \sum_{j=0}^{n-1} \mathbb{E} \left\| \phi_j(W_{E_{t_{j+1}}} - W_{E_{t_j}}) \right\|_H^2
$$

$$
+ \sum_{i \neq j=0}^{n-1} \mathbb{E} \left\langle \phi_j(W_{E_{t_{j+1}}} - W_{E_{t_j}}), \phi_i(W_{E_{t_{i+1}}} - W_{E_{t_i}}) \right\rangle_H := I + II.
$$
Let \( \{e_m\}_{m=1}^{\infty} \) in \( H \) and \( \{f_l\}_{l=1}^{\infty} \) in \( K \) be ONBs. For fixed \( j \), \( I_j \) is denoted by

\[
I_j = \mathbb{E} \left\| \phi_j(W_{E_{t,j+1}} - W_{E_{t,j}}) \right\|_H^2 = \mathbb{E} \left\langle \phi_j(W_{E_{t,j+1}} - W_{E_{t,j}}), e_m \right\rangle_H^2
\]

\[
= \sum_{m=1}^{\infty} \mathbb{E} \left( \mathbb{E} \left( \left\langle \phi_j(W_{E_{t,j+1}} - W_{E_{t,j}}), e_m \right\rangle_H^2 | G_{t_j} \right) \right)
\]

\[
= \sum_{m=1}^{\infty} \mathbb{E} \left( \mathbb{E} \left( \left\langle W_{E_{t,j+1}} - W_{E_{t,j}}, \phi_j^* e_m \right\rangle_K^2 | G_{t_j} \right) \right)
\]

\[
= \sum_{m=1}^{\infty} \mathbb{E} \left( \mathbb{E} \left( \left( \sum_{l=1}^{\infty} \left\langle W_{E_{t,j+1}} - W_{E_{t,j}}, f_l \right\rangle_K \left\langle \phi_j^* e_m, f_l \right\rangle_K^2 \right) | G_{t_j} \right) \right)
\]

\[
+ \sum_{m=1}^{\infty} \mathbb{E} \left( \mathbb{E} \left( \left( \sum_{l \neq l'}^{\infty} \left\langle W_{E_{t,j+1}} - W_{E_{t,j}}, f_l \right\rangle_K \langle \phi_j^* e_m, f_{l'} \rangle_K \right) \right) G_{t_j} \right)
\]

\[
= J_1 + J_2.
\]
Since \( \phi_j^* \) is \( G_{t_j} \)-measurable and \( W_{E_{t_j+1}} \) is a discrete martingale with respect to \( G_{t_j} \), the first term, \( J_1 \), becomes

\[
J_1 = \sum_{m=1}^{\infty} \mathbb{E}\left( \sum_{l=1}^{\infty} \langle \phi_j^* e_m, f_l \rangle_K^2 \mathbb{E}\left( \left\langle W_{E_{t_j+1}} - W_{E_{t_j}}, f_l \right\rangle_K^2 \big| G_{t_j} \right) \right)
= \sum_{m=1}^{\infty} \mathbb{E}\left\{ \sum_{l=1}^{\infty} \langle \phi_j^* e_m, f_l \rangle_K^2 \mathbb{E}\left( \left\langle W_{E_{t_j+1}}, f_l \right\rangle_K^2 \big| G_{t_j} \right) \right\}
= \sum_{m=1}^{\infty} \mathbb{E}\left\{ \sum_{l=1}^{\infty} \langle \phi_j^* e_m, f_l \rangle_K \langle W_{E_{t_j}}, f_l \rangle_K \right\}

- \sum_{m=1}^{\infty} \mathbb{E}\left\{ \sum_{l=1}^{\infty} \langle \phi_j^* e_m, f_l \rangle_K \langle W_{E_{t_j+1}}, f_l \rangle_K \right\}

= \mathbb{E}\left\{ \sum_{m=1}^{\infty} \left( E_{t_{j+1}} - E_{t_j} \right) \sum_{l=1}^{\infty} \lambda_l \langle \phi_j^* e_m, f_l \rangle_K \right\}

= \mathbb{E}\left\{ \left( E_{t_{j+1}} - E_{t_j} \right) \sum_{m,l=1}^{\infty} \langle \phi_j (\lambda_l^{1/2} f_l), e_m \rangle_H^2 \right\}

= \mathbb{E}\left\{ \left( E_{t_{j+1}} - E_{t_j} \right) \| \phi_j \|_{L_2(K_Q,H)}^2 \right\}.

Also using the \( G_{t_j} \)-measurability of \( \phi_j^* \) and Lemma 3.3, the second term, \( J_2 \), becomes

\[
J_2 = \sum_{m=1}^{\infty} \mathbb{E}\left\{ \sum_{l \neq l'}^{\infty} \langle \phi_j^* e_m, f_l \rangle_K \langle \phi_j^* e_m, f_{l'} \rangle_K \right\}

\times \mathbb{E}\left\{ \left\langle W_{E_{t_{j+1}}}, f_l \right\rangle_K \left\langle W_{E_{t_{j+1}}}, f_{l'} \right\rangle_K \big| G_{t_j} \right\} = 0.

Thus,

\[
I = \sum_{j=0}^{n-1} I_j = \sum_{j=0}^{n-1} \mathbb{E}\left\{ \left( E_{t_{j+1}} - E_{t_j} \right) \| \phi_j \|_{L_2(K_Q,H)}^2 \right\}

= \mathbb{E} \int_0^T \| \Phi(s) \|_{L_2(K_Q,H)}^2 dE_s < \infty.

On the other hand, without loss of generality, assume that \( i < j \).

From Lemma 3.3,

\[
II = \mathbb{E} \sum_{m=1}^{\infty} \mathbb{E} \left( \sum_{l,l'=1}^{\infty} \left< W_{E_{l'}+1} - W_{E_l}, f_l \right> \phi_j^* e_m, f_l \right)_K \\
\times \left< W_{E_{l'+1}} - W_{E_{l'}}, f_l' \right> \phi_i^* e_m, f_l' \right)_K | G_{t_j} = 0.
\]

\[\square\]

**Theorem 3.5.** (Time-changed Itô Isometry) For \( t \in [0,T] \), the stochastic integral \( \Phi \mapsto \int_0^t \Phi(s)dW_E \) with respect to a \( K \)-valued time-changed \( Q \)-Wiener process \( W_{E_t} \) is an isometry between \( \tilde{\Lambda}_2(K_Q,H) \) and the space of continuous square-integrable martingales \( M^2_T(H) \), i.e.,

\[
\mathbb{E} \left\| \int_0^t \Phi(s)dW_E \right\|^2_H = \mathbb{E} \int_0^t \left\| \Phi(s) \right\|^2_{\mathcal{L}_2(K_Q,H)}dE_s < \infty. \tag{3.2}
\]

**Proof.** For elementary processes \( \Phi \in \mathcal{E}(\mathcal{L}(K,H)) \), Theorem 3.4 establishes the desired equality (3.2) and consequently the square-integrability of the integral \( \int_0^t \Phi(s)dW_E \). Furthermore, since the time-changed \( Q \)-Wiener process, \( W_{E_t} \), is a \( K \)-valued martingale, for any \( h \in H \) and \( s < t \),

\[
\mathbb{E} \left( \left< \int_0^t \Phi(r)dW_{E_r}, h \right>_H \big| G_s \right) = \mathbb{E} \left( \left< \sum_{j=0}^{n-1} \phi_j(W_{E_{t_j+1}\wedge t} - W_{E_{t_j}\wedge t}), h \right>_H \big| G_s \right) \\
= \sum_{j=0}^{n-1} \mathbb{E} \left( \left< W_{E_{t_j+1}\wedge t} - W_{E_{t_j}\wedge t}, \phi_j^*(h) \right> _K \big| G_s \right) \\
= \sum_{j=0}^{n-1} \left< W_{E_{t_j+1}\wedge s} - W_{E_{t_j}\wedge s}, \phi_j^*(h) \right> _K \\
= \left< \sum_{j=0}^{n-1} \phi_j(W_{E_{t_j+1}\wedge s} - W_{E_{t_j}\wedge s}), h \right> _H = \left< \int_0^s \Phi(s)dW_E, h \right>_H.
\]
which implies that the stochastic integral $\int_0^t \Phi(s) dW_E_s$ is a square-integrable martingale. Therefore, the desired result holds when $\Phi(s)$ is an elementary process.

Now, let $\{\Phi_n\}_{n=1}^\infty$ be a sequence of elementary processes approximating $\Phi \in \tilde{\Lambda}_2(K_Q, H)$. Assume that $\Phi_1 = 0$ and $||\Phi_{n+1} - \Phi_n||_{\tilde{\Lambda}_2(K_Q, H)} < \frac{1}{2^n}$. Then, by Doob’s Maximal Inequality,

$$
\sum_{n=1}^\infty \mathbb{P} \left( \sup_{t \leq T} \left\| \int_0^t \Phi_{n+1}(s) dW_E_s - \int_0^t \Phi_n(s) dW_E_s \right\|_H > \frac{1}{n^2} \right) \leq \sum_{n=1}^\infty n^4 \mathbb{E} \left( \int_0^T (\Phi_{n+1}(s) - \Phi_n(s)) dW_E_s \right)^2 
$$

$$
= \sum_{n=1}^\infty n^4 \mathbb{E} \int_0^T ||\Phi_{n+1}(s) - \Phi_n(s)||_{\mathcal{L}(K_Q, H)}^2 dE_s \leq \frac{T^\beta}{\Gamma(\beta + 1)} \sum_{n=1}^\infty \frac{n^4}{2^n} < \infty.
$$

By Borel-Cantelli, it follows that for some $k(\omega) > 0$,

$$
\sup_{t \leq T} \left\| \int_0^t \Phi_{n+1}(s) dW_E_s - \int_0^t \Phi_n(s) dW_E_s \right\|_H \leq \frac{1}{n^2}, \ n > k(\omega),
$$

holds $\mathbb{P}$-almost surely. Therefore, for every $t \leq T$,

$$
\sum_{n=1}^\infty \left( \int_0^t \Phi_{n+1}(s) dW_E_s - \int_0^t \Phi_n(s) dW_E_s \right) \to \int_0^t \Phi(s) dW_E_s \text{ in } L^2(\Omega, H),
$$

which also converges $\mathbb{P}$-a.s. to a continuous version of the integral.

Thus, the map $\Phi \mapsto \int_0^t \Phi(s) dW_E_s$, viewed as an isometry from elementary processes to the space of continuous square-integrable martingales, has an extension to $\Phi \in \tilde{\Lambda}_2(K_Q, H)$ by the completeness property of $H$. \hfill \Box

The following two change of variable formulas concern the Itô stochastic integral related to the time-change $E_t$. They are needed
later and can be considered as the Hilbert space extensions of formulas in \[26\].

**Theorem 3.6.** \((1st \ change \ of \ variable \ formula \ in \ Hilbert \ space)\) Let \(W_t\) be a \(Q\)-Wiener process in a separable Hilbert space \(K\), \(\Phi \in \tilde{\Lambda}_2(K_Q, H)\), and \(E_t\) be the inverse of a \(\beta\)-stable subordinator. Then, with probability one, for all \(t \geq 0\),

\[
\int_0^{E_t} \Phi(s) dW_s = \int_0^t \Phi(E_s) dW_{E_s}.
\]

**Proof.** Let \(\{f_j\}_{j=1}^\infty\) be an ONB in the separable Hilbert space \(K\) consisting of eigenvectors of \(Q\). Then, it follows from Lemma 3.1 that

\[
\int_0^{E_t} \Phi(s) dW_s = \sum_{j=1}^\infty \int_0^{E_t} (\Phi(s)\lambda_j^{1/2} f_j) d\langle W_s, \lambda_j^{1/2} f_j \rangle_{K_Q}.
\]

For any \(h \in H\),

\[
\left< \int_0^{E_t} \Phi(s) dW_s, h \right>_H = \left< \sum_{j=1}^\infty \int_0^{E_t} (\Phi(s)\lambda_j^{1/2} f_j) d\langle W_s, \lambda_j^{1/2} f_j \rangle_{K_Q}, h \right>_H
\]

\[
= \sum_{j=1}^\infty \int_0^{E_t} \langle (\Phi(s)\lambda_j^{1/2} f_j), h \rangle_H d\langle W_s, \lambda_j^{1/2} f_j \rangle_{K_Q}
\]

\[
= \sum_{j=1}^\infty \int_0^{E_t} \langle (E_s)\lambda_j^{1/2} f_j), h \rangle_H d\langle W_{E_s}, \lambda_j^{1/2} f_j \rangle_{K_Q}
\]

\[
= \sum_{j=1}^\infty \int_0^{E_t} \langle (E_s)\lambda_j^{1/2} f_j), h \rangle_H d\langle W_{E_s}, \lambda_j^{1/2} f_j \rangle_{K_Q}
\]

\[
= \left< \sum_{j=1}^\infty \int_0^{E_t} (\Phi(s)\lambda_j^{1/2} f_j) d\langle W_{E_s}, \lambda_j^{1/2} f_j \rangle_{K_Q}, h \right>_H
\]

\[
= \left< \int_0^{E_t} \Phi(s) dW_s, h \right>_H.
\]

The third equality follows from the first change of variable formula of the real-valued stochastic integral from \[26\]. \(\square\)
Theorem 3.7. (2nd change of variable formula in Hilbert space)

Let \( W_t \) be a Q-Wiener process in a separable Hilbert space \( K \) and \( \Phi \in \tilde{\Lambda}_2(K_Q,H) \). Let \( U_t \) be a \( \beta \)-stable subordinator with \( \beta \in (0,1) \) and \( E_t \) be its inverse stable subordinator. Then, with probability one, for all \( t \geq 0 \),

\[
\int_0^t \Phi(s) dW_{E_s} = \int_0^{E_t} \Phi(U_{s-}) dW_s.
\]

Proof. Let \( \{f_j\}_{j=1}^\infty \) be an ONB in the separable Hilbert space \( K \) consisting of eigenvectors of \( Q \). Applying Definition 3.2 yields,

\[
\int_0^t \Phi(s) dW_s = \sum_{j=1}^\infty \int_0^t (\Phi(s)\lambda_j^{1/2}f_j)d\langle W_{E_s},\lambda_j^{1/2}f_j\rangle_{K_Q}.
\]

For any \( h \in H \),

\[
\left\langle \int_0^t \Phi(s) dW_{E_s}, h \right\rangle_H = \left\langle \sum_{j=1}^\infty \int_0^t (\Phi(s)\lambda_j^{1/2}f_j)d\langle W_{E_s},\lambda_j^{1/2}f_j\rangle_{K_Q}, h \right\rangle_H
\]

\[
= \sum_{j=1}^\infty \int_0^t \left\langle (\Phi(s)\lambda_j^{1/2}f_j), h \right\rangle_H d\langle W_{E_s},\lambda_j^{1/2}f_j\rangle_{K_Q}
\]

\[
= \sum_{j=1}^\infty \int_0^{E_t} \left\langle (\Phi(U_{s-})\lambda_j^{1/2}f_j), h \right\rangle_H d\langle W_s,\lambda_j^{1/2}f_j\rangle_{K_Q}
\]

\[
= \left\langle \sum_{j=1}^\infty \int_0^{E_t} (\Phi(U_{s-})\lambda_j^{1/2}f_j)d\langle W_s,\lambda_j^{1/2}f_j\rangle_{K_Q}, h \right\rangle_H
\]

\[
= \left\langle \int_0^{E_t} \Phi(U_{s-}) dW_s, h \right\rangle_H.
\]

The first equality follows from Definition 3.2 and the third equality follows from the second change of variable formula for real-valued stochastic integrals from [26]. \( \square \)
3.2. Time-changed Itô formula in Hilbert space

The technique used to develop the time-changed Itô formula in this section is inspired by the proof of the standard Itô formula of Theorem 2.9 in [5].

**Theorem 3.8.** (Time-changed Itô formula) Let $Q$ be a symmetric, nonnegative definite trace-class operator on a separable Hilbert space $K$, and let $\{W_{E_t}\}_{0 \leq t \leq T}$ be a time-changed $Q$-Wiener process on a filtered probability space $(\Omega, \mathcal{G}, \{\mathcal{G}_t\}_{0 \leq t \leq T}, \mathbb{P})$. Assume that a stochastic process $X(t)$ is given by

$$X(t) = X(0) + \int_0^t \psi(s)ds + \int_0^t \gamma(s)dE_s + \int_0^t \phi(s)dW_{E_s},$$

where $X(0)$ is a $\mathcal{G}_0$-measurable $H$-valued random variable, $\psi(s)$ and $\gamma(s)$ are $H$-valued $\mathcal{G}_s$-measurable $\mathbb{P}$-a.s. integrable processes on $[0, T]$ such that

$$\int_0^T \|\psi(s)\|_H ds < \infty \text{ and } \int_0^T \|\gamma(s)\|_H dE_s < \infty,$$

and $\phi \in \hat{L}_2(K_Q, H)$. Also assume that $F : H \to \mathbb{R}$ is continuous and its Fréchet derivatives $F_x : H \to \mathcal{L}(H, \mathbb{R})$ and $F_{xx} : H \to \mathcal{L}(H, \mathcal{L}(H, \mathbb{R}))$ are continuous and bounded on bounded subsets of $H$. 
Then,

\[
F(X(t)) - F(X(0)) = \int_0^t \left\langle F_x(X(s)), \psi(s) \right\rangle_H ds \\
+ \int_0^{E_t} \left\langle F_x(X(U(s-))), \gamma(U(s-)) \right\rangle_H ds \\
+ \int_0^{E_t} \left\langle F_x(X(U(s-))), \phi(U(s-))dW_s \right\rangle_H \\
+ \frac{1}{2} \int_0^{E_t} \text{tr}(F_{xx}(X(U(s-)))(\phi(U(s-))Q^{1/2})(\phi(U(s-))Q^{1/2})^*)ds,
\]

\(\mathbb{P}\)-a.s. for all \(t \in [0,T]\).

**Proof.** First, the desired Itô formula is reduced to the case where \(\psi(s) = \psi\), \(\gamma(s) = \gamma\), and \(\phi(s) = \phi\) are constant processes for \(s \in [0,T]\). Let \(C > 0\) be a constant, and define the stopping time

\[
\tau_C = \inf \left\{ t \in [0,T] : \max \left( \|X(t)\|_H, \int_0^t \|\phi(s)\|_H ds, \int_0^t \|\gamma(s)\|_H dE_s, \int_0^t \|\phi(s)\|_{L_2(K_Q,H)}^2 dE_s \right) \geq C \right\}.
\]

Then, define \(X_C(t)\) as

\[
X_C(t) = X_C(0) + \int_0^t \psi_C(s) ds + \int_0^t \gamma_C(s) dE_s + \int_0^t \phi_C(s) dW_{E_s}, \quad t \in [0,T],
\]

where \(X_C(t) = X(t \wedge \tau_C), \psi_C(t) = \psi(t)1_{[0,\tau_C]}(t), \gamma_C(t) = \gamma(t)1_{[0,\tau_C]}(t), \) and \(\phi_C(t) = \phi(t)1_{[0,\tau_C]}(t)\). It is enough to prove the Itô formula for the processes stopped at \(\tau_C\). Since

\[
\mathbb{P}\left( \int_0^T \|\psi_C(s)\|_H ds < \infty \right) = 1, \quad \mathbb{P}\left( \int_0^T \|\gamma_C(s)\|_H dE_s < \infty \right) = 1, \quad \text{and} \quad E \int_0^T \|\phi_C(s)\|_{L_2(K_Q,H)}^2 dE_s < \infty,
\]
it follows that $\psi_C$, $\gamma_C$, and $\phi_C$ can be approximated respectively by sequences of bounded elementary processes $\psi_{C,n}$, $\gamma_{C,n}$, and $\phi_{C,n}$ such that as $n \to \infty$

$$
\int_0^t \| \psi_{C,n}(s) - \psi_C(s) \|_H ds \to 0, \quad \int_0^t \| \gamma_{C,n}(s) - \gamma_C(s) \|_H dE_s \to 0, \quad \text{and}
$$

$$
\left\| \int_0^t \phi_{C,n}(s)dW_{E_s} - \int_0^t \phi_C(s)dW_{E_s} \right\|_H \to 0, \quad P \text{- a.s.}
$$

uniformly in $t \in [0, T]$. Let

$$
X_{C,n}(t) = X(0) + \int_0^t \psi_{C,n}(s)ds + \int_0^t \gamma_{C,n}(s)dE_s + \int_0^t \phi_{C,n}(s)dW_{E_s}.
$$

Then, as $n \to \infty$

$$
\sup_{t \leq T} \| X_{C,n}(t) - X_C(t) \|_H \to 0, \quad P \text{- a.s.} \quad (3.4)
$$

Assume the Itô formula for $X_{C,n}(t)$ holds $P$-a.s. for all $t \in [0, T]$, i.e.,

$$
F(X_{C,n}(t)) - F(X(0)) = \int_0^t \langle F_x(X_{C,n}(s)), \phi_{C,n}(s)dW_{E_s} \rangle_H + \int_0^t \langle F_x(X_{C,n}(s)), \psi_{C,n}(s) \rangle_H ds + \int_0^t \langle F_x(X_{C,n}(s)), \gamma_{C,n}(s) \rangle_H dE_s
$$

$$
+ \int_0^t \frac{1}{2} \text{tr}[F_{xx}(X_{C,n}(s)) (\phi_{C,n}(s)Q^{1/2})(\phi_{C,n}(s)Q^{1/2})^*] dE_s
$$

$$
:= I_{C,n}^1 + I_{C,n}^2 + I_{C,n}^3 + I_{C,n}^4.
$$

By using the continuity of $F$ and the continuity and local boundedness of $F_x$ and $F_{xx}$, it will suffice to show that the following holds
\[ F(X_{C}(t)) - F(X(0)) = \int_{0}^{t} \langle F_{x}(X_{C}(s)), \phi_{C}(s) \rangle dW_{E_{s}} \]
\[+ \int_{0}^{t} \langle F_{x}(X_{C}(s)), \psi_{C}(s) \rangle H \, ds + \int_{0}^{t} \langle F_{x}(X_{C}(s)), \gamma_{C}(s) \rangle H \, dE_{s} \]
\[+ \frac{1}{2} \int_{0}^{t} \text{tr}[F_{xx}(X_{C}(s))(\phi_{C}(s)Q^{1/2})(\phi_{C}(s)Q^{1/2})^{*}] dE_{s} \]
\[= I_{1}^{C} + I_{2}^{C} + I_{3}^{C} + I_{4}^{C}. \] (3.6)

Consider, term by term, the difference between both sides of (3.5) and (3.6). Due to the continuity of \( F \) and almost sure convergence in (3.4), the left hand side of (3.5) converges to the left hand side of (3.6) \( \mathbb{P} \)-a.s for all \( t \leq T \), i.e.,
\[ F(X_{C,n}(t)) \to F(X_{C}(t)), \quad \mathbb{P} - \text{a.s. as } n \to \infty. \] (3.7)

Turn to the first terms in both right hand sides of (3.5) and (3.6),
\[ \mathbb{E}|I_{1}^{C,n} - I_{1}^{C}|^{2} = \mathbb{E} \left| \int_{0}^{t} \left( \phi_{C,n}^{*}(s)F_{x}(X_{C,n}(s)) - \phi_{C}^{*}(s)F_{x}(X_{C}(s)) \right) dW_{E_{s}} \right|^{2} \]
\[\leq 2 \mathbb{E} \int_{0}^{t} \| (\phi_{C,n}^{*}(s) - \phi_{C}^{*}(s))F_{x}(X_{C,n}(s)) \|_{L_{2}(K_{Q},R)}^{2} dE_{s} \]
\[+ 2 \mathbb{E} \int_{0}^{t} \| \phi_{C}^{*}(s)(F_{x}(X_{C,n}(s)) - F_{x}(X_{C}(s))) \|_{L_{2}(K_{Q},R)}^{2} dE_{s} \]
\[\leq 2 \mathbb{E} \int_{0}^{t} \left( \| \phi_{C,n}^{*}(s) \|_{L_{2}(K_{Q},H)}^{2} \| F_{x}(X_{C,n}(s)) \|_{H}^{2} \right) dE_{s} \]
\[+ 2 \mathbb{E} \int_{0}^{t} \left( \| \phi_{C}^{*}(s) \|_{L_{2}(K_{Q},H)}^{2} \| F_{x}(X_{C,n}(s)) - F_{x}(X_{C}(s)) \|_{H}^{2} \right) dE_{s} \]
\[:= J_{1} + J_{2}, \]

where \( \int_{0}^{t} \beta^{*}(s)\alpha(s) dW_{E_{s}} := \int_{0}^{t} \langle \alpha(s), \beta(s) \rangle dW_{E_{s}} \) and \( \beta^{*}(s) \) is the adjoint operator of \( \beta(s) \). Since \( F_{x} \) is bounded on bounded subsets of
$H$, there exists an $M > 0$ such that

$$J_1 \leq ME \int_0^t \| \phi_C^*(s) - \phi_{C,n}^* \|^2_{L^2(K_Q,H)} dE_s \to 0, \text{ as } n \to \infty.$$ 

Since $\phi_C(s)$ is square integrable in the space $\tilde{\Lambda}_2(K_Q,H)$ and $F_x$ is bounded in $H$, $J_2 \to 0$ by applying the Lebesgue dominated convergence theorem. So, $I_{C,n}^1$ converges to $I_C^1$ in mean square, i.e.,

$$\mathbb{E}|I_{C,n}^1 - I_C^1|^2 \to 0, \quad (3.8)$$

and thus converges in probability. For the second terms, $I_{C,n}^2$ and $I_C^2$, the RHSs of (3.5) and (3.6), applying the conditions of (3.3) and (3.4) leads to

$$I_{C,n}^2 - I_C^2 = \int_0^t \left( \left\langle F_x(X_{C,n}(s)) - F_x(X_C(s)), \psi_{C,n}(s) \right\rangle_{H} 
+ \left\langle F_x(X_C(s)), \psi_{C,n}(s) - \psi_C(s) \right\rangle_{H} \right) ds \to 0, \mathbb{P} - \text{a.s.} \quad (3.9)$$

Similarly, for the third terms, $I_{C,n}^3$ and $I_C^3$, the RHSs of (3.5) and (3.6),

$$I_{C,n}^3 - I_C^3 = \int_0^t \left( \left\langle F_x(X_{C,n}(s)) - F_x(X_C(s)), \gamma_{C,n}(s) \right\rangle_{H} 
+ \left\langle F_x(X_C(s)), \gamma_{C,n}(s) - \gamma_C(s) \right\rangle_{H} \right) dE_s \to 0, \mathbb{P} - \text{a.s.} \quad (3.10)$$

Before proceeding to the fourth terms, $I_{C,n}^4$ and $I_C^4$, of RHSs of (3.5) and (3.6), note that

$$\| \phi_{C,n}(s) - \phi_C(s) \|_{\tilde{\Lambda}_2(K_Q,H)} \to 0,$$
which means there exists a subsequence $n_k$ such that for all $s \leq T$

$$
\|\phi_{C,n_k}(s) - \phi_C(s)\|_{L_2(K_Q,H)} \to 0, \mathbb{P} - \text{a.s.}
$$

Thus, for the eigenvectors $\{f_j\}_{j=1}^\infty$ of $Q$ and all $t \leq T$,

$$
\|\phi_{C,n_k}(s)f_j - \phi_C(s)f_j\|_H \to 0, \mathbb{P} - \text{a.s.} \quad (3.11)
$$

On the other hand, for the ONB, $\{f_j\}_{j=1}^\infty$, in the Hilbert space $K$,

$$
tr(F_{xx}(X_{C,n_k}(s))\phi_{C,n_k}(s)Q\phi_{C,n_k}^*(s)) = tr(\phi_{C,n_k}^*(s)F_{xx}(X_{C,n_k}(s))\phi_{C,n_k}(s)Q)
$$

$$
= \sum_{j=1}^\infty \lambda_j \langle F_{xx}(X_{C,n_k}(s))\phi_{C,n_k}(s)f_j, \phi_{C,n_k}(s)f_j \rangle_H,
$$

where $\lambda_j$ is the eigenvalue associated with eigenvector $f_j$ of $Q$. Since $X_{C,n_k}(s)$ is bounded and $F_{xx}$ is continuous, (3.11) implies that for $s \leq T$

$$
\langle F_{xx}(X_{C,n_k}(s))\phi_{C,n_k}(s)f_j, \phi_{C,n_k}(s)f_j \rangle_H
$$

$$
\to \langle F_{xx}(X_C(s))\phi_C(s)f_j, \phi_C(s)f_j \rangle_H, \mathbb{P} - \text{a.s.}
$$

By the Lebesgue Dominated Convergence Theorem (DCT) (with respect to the counting measure), it holds a.e. on $[0,T] \times \Omega$ that

$$
tr(F_{xx}(X_{C,n_k}(s))\phi_{C,n_k}(s)Q\phi_{C,n_k}^*(s))
$$

$$
\to tr(F_{xx}(X_C(s))\phi_C(s)Q\phi_C^*(s)). \quad (3.12)
$$

Moreover, the left hand side of (3.12) is bounded above by

$$
\eta_n(s) := \|F_{xx}(X_{C,n_k}(s))\|_{L(H)}\|\phi_{C,n_k}\|_{L_2(K_Q,H)}^2;
$$
and a.e. on $[0, T] \times \Omega$

$$\eta_n(s) \to \eta(s) = \|F_{xx}(X_C(s))\|_{L(H)} \|\phi_C\|_{L_2(K_Q,H)}^2.$$ 

So, by the boundedness of $F_{xx}$, $\int_0^t \eta_n(s) dE_s \to \int_0^t \eta(s) dE_s$ and by the general Lebesgue DCT, it holds $\mathbb{P}$-a.s. that for $t \leq T$

$$I_{C,n_k}^4 - I_C^4 = \int_0^t \frac{1}{2} tr[F_{xx}(X_{C,n_k}(s))(\phi_{C,n_k}(s)Q^{1/2})(\phi_{C,n_k}(s)Q^{1/2})^*]dE_s - \int_0^t \frac{1}{2} tr[F_{xx}(X_C(s))(\phi_{C}(s)Q^{1/2})(\phi_{C}(s)Q^{1/2})^*]dE_s \to 0. \tag{3.13}$$

Therefore, from (3.7), (3.8), (3.9), (3.10) and (3.13), the Itô formula for the process $X_{C,n}(t)$, (3.5), converges in probability to the Itô formula for the process $X_C(t)$, (3.6), and possibly for a subsequence, $n_k$, converges $\mathbb{P}$-a.s.

Second, the proof can be reduced to the case where

$$X(t) = X(0) + \psi t + \gamma E_t + \phi W_{E_t}$$

where $\psi, \gamma,$ and $\phi$ are $\mathcal{G}_0$-measurable random variables independent of $t$. Define the function $u(t_1, t_2, x) : \mathbb{R}_+ \times \mathbb{R}_+ \times H \to \mathbb{R}$ as

$$u(t, E_t, W_{E_t}) = F(X(0) + \psi t + \gamma E_t + \phi W_{E_t}) = F(X(t)).$$

Now, we prove that the Itô formula holds for the function $u(t_1, t_2, x)$. First, let $0 = t_1 < t_2 < \ldots < t_n = t \leq T$ be a partition of an interval
If \( [0, t] \), then

\[
\begin{align*}
    u(t, E_t, W_{E_t}) - u(0, 0, 0) &= \sum_{j=1}^{n-1} [u(t_{j+1}, E_{t_{j+1}}, W_{E_{t_{j+1}}}) - u(t_j, E_{t_j}, W_{E_{t_j}})] \\
    &\quad + \sum_{j=1}^{n-1} [u(t_j, E_{t_j}, W_{E_{t_j}}) - u(t_j, E_{t_j}, W_{E_{t_j}})] \\
    &\quad + \sum_{j=1}^{n-1} [u(t_j, E_{t_j}, W_{E_{t_j}}) - u(t_j, E_{t_j}, W_{E_{t_j}})].
\end{align*}
\]

Also, let \( \Delta t_j = t_{j+1} - t_j \), \( \Delta E_j = E_{t_{j+1}} - E_{t_j} \) and \( \Delta W_j = W_{E_{t_{j+1}}} - W_{E_{t_j}} \). Let \( \theta_j \in [0, 1] \) be a random variable, and \( \tilde{t}_j = t_j + \theta_j(t_{j+1} - t_j) \), \( \tilde{E}_j = E_{t_j} + \theta_j(E_{t_{j+1}} - E_{t_j}) \) and \( \tilde{W}_j = W_{E_{t_j}} + \theta_j(W_{E_{t_{j+1}}} - W_{E_{t_j}}) \).
Using Taylor’s formula,

\[ u(t, E_t, W_{E_t}) - u(0, 0, 0) \]

\[ = \sum_{j=1}^{n-1} u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) \Delta t_j + \sum_{j=1}^{n-1} u_{t_2}(t_j, \bar{E}_j, W_{E_{t_{j+1}}}) \Delta E_j \]

\[ + \sum_{j=1}^{n-1} [\langle u_x(t_j, E_{t_j}, W_{E_{t_j}}), \Delta W_j \rangle_K + \frac{1}{2} \langle u_{xx}(t_j, E_{t_j}, \bar{W}_j)(\Delta W_j), \Delta W_j \rangle_K] \]

\[ = \sum_{j=1}^{n-1} u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) \Delta t_j \]

\[ + \sum_{j=1}^{n-1} u_{t_2}(t_j, E_{t_j}, W_{E_{t_{j+1}}}) \Delta E_j + \sum_{j=1}^{n-1} \langle u_x(t_j, E_{t_j}, W_{E_{t_j}}), \Delta W_j \rangle_K \]

\[ + \frac{1}{2} \sum_{j=1}^{n-1} \langle u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j), \Delta W_j \rangle_K \]

\[ + \sum_{j=1}^{n-1} [u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) - u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}})] \Delta t_j \]

\[ + \sum_{j=1}^{n-1} [u_{t_2}(t_j, \bar{E}_{j+1}, W_{E_{t_{j+1}}}) - u_{t_2}(t_j, E_{t_j}, W_{E_{t_{j+1}}})] \Delta E_j \]

\[ + \frac{1}{2} \sum_{j=1}^{n-1} [\langle u_{xx}(t_j, E_{t_j}, \bar{W}_j)(\Delta W_j) - u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j), \Delta W_j \rangle_K] \]

\[ := I_1 + I_2 + I_3 + I_4 + I_5 + I_6 + I_7. \]

By the uniform continuity of the mappings:

\[ [0, T] \times [0, T] \times [0, T] \ni (t, s, r) \rightarrow u_{t_1}(t, E_s, W_{E_r}) \in \mathbb{R} \]

\[ [0, T] \times [0, T] \times [0, T] \ni (t, s, r) \rightarrow u_{t_2}(t, E_s, W_{E_r}) \in \mathbb{R} \]
and the continuity of the map \([0, T] \ni t \to u_x(t, E_t, W_{E_t}) \in K^*\), the following holds \(\mathbb{P}\text{-a.s.}\):

\[
I_1 = \sum_{j=1}^{n-1} u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) \Delta t_j \to \int_0^t u_{t_1}(s, E_s, W_{E_s}) ds,
\]

\[
I_2 = \sum_{j=1}^{n-1} u_{t_2}(t_j, E_{t_{j+1}}) \Delta E_{t_j} \to \int_0^t u_{t_2}(s, E_s, W_{E_s}) dE_s, \quad (3.15)
\]

\[
I_3 = \sum_{j=1}^{n-1} \langle u_x(s, E_{E_s}), \Delta W_j \rangle_K \to \int_0^t \langle u_x(s, E_s), dW_{E_s} \rangle_K.
\]

Also since the time change \(E_t\) has bounded variation,

\[
|I_5| = \left| \sum_{j=1}^{n-1} [u_{t_1}(\bar{t}_j, E_{\bar{t}_{j+1}}, W_{E_{\bar{t}_{j+1}}}) - u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}})] \Delta t_j \right|
\]

\[
\leq T \sup_{j \leq n-1} \left| u_{t_1}(\bar{t}_j, E_{\bar{t}_{j+1}}, W_{E_{\bar{t}_{j+1}}}) - u_{t_1}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) \right| \to 0,
\]

\[
|I_6| = \left| \sum_{j=1}^{n-1} [u_{t_2}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) - u_{t_2}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}})] \Delta E_j \right| \quad (3.16)
\]

\[
\leq \sum_{j=1}^{n-1} |\Delta E_{t_j}| \sup_{j \leq n-1} \left| u_{t_2}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) - u_{t_2}(t_j, E_{t_{j+1}}, W_{E_{t_{j+1}}}) \right| \to 0.
\]
Similarly, by the continuity of the map $K \ni x \to u_{xx}(t_1, t_2, x) \in \mathcal{L}(K, K)$,

$$
|I_7| = \frac{1}{2} \sum_{j=1}^{n-1} \left| [u_{xx}(t_j, E_{t_j}, W_j)(\Delta W_j) - u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j)]_K \right|
$$

$$
\leq \sup_{j \leq n-1} \|u_{xx}(t_j, E_{t_j}, \bar{W}_j)(\Delta W_j) - u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j)\|_{\mathcal{L}(K, K)} \times \sum_{j=1}^{n-1} \|\Delta W_j\|^2_K \to 0
$$

with probability one as $n \to \infty$ since the function $u$ has the same smoothness as $F$ and $\bar{W}_j \to W_{E_{t_j}}$ as the increments $t_{j+1} - t_j$ get smaller. It remains to deal with the fourth term, $I_4$. Let $1^N_j = 1_{\{\max\{\|W_{E_{t_j}}\|_K \leq N, i \leq j\}\}}$ which is $\mathcal{G}_{t_j}$-measurable. To handle $I_4$, the following computations are helpful. First,

$$
\mathbb{E}\left( 1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j), \Delta W_j)_K \left| \mathcal{G}_{t_j} \right) = \mathbb{E}\left( 1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) \sum_{k=1}^{\infty} \lambda_k^{1/2} (w_k(E_{t_{j+1}}) - w_k(E_{t_j})) f_k, \right. \\
\sum_{l=1}^{\infty} \lambda_l^{1/2} (w_l(E_{t_{j+1}}) - w_l(E_{t_j})) f_l \left| \mathcal{G}_{t_j} \right) \\
= \sum_{k=1}^{\infty} \mathbb{E}\left( \lambda_k (1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) f_k, f_k)_K (w_k(E_{t_{j+1}}) - w_k(E_{t_j}))^2 \left| \mathcal{G}_{t_j} \right) \\
= tr(1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) \mathcal{Q}) \Delta E_j.
$$
Second, for the cross term arising in the computation below, without loss of generality, assume $i < j$. Then,

\[
I^N := \mathbb{E}\left\{ \left( \langle 1^N_i u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_i), \Delta W_i \rangle_K 
- \text{tr}\left(1^N_i u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q\right)\Delta E_i \right) \times \left( \langle 1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j), \Delta W_j \rangle_K 
- \text{tr}\left(1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q\right)\Delta E_j \right) \right\} \\
= \mathbb{E}\left\{ \left( \langle 1^N_i u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_i), \Delta W_i \rangle_K 
- \text{tr}\left(1^N_i u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q\right)\Delta E_i \right) \times \mathbb{E}\left( \langle 1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})(\Delta W_j), \Delta W_j \rangle_K 
- \text{tr}\left(1^N_j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q\right)\Delta E_j \big| G_{t(i+1)} \right) \right\}.
\]
Applying the tower property of conditional expectation and (3.18) yields

\[ I^N = \mathbb{E} \left\{ \left( \langle 1_i^N u_{xx}(t_j, E_{t_j}, W_{E_{t_i}}) \delta W_i, \delta W_i \rangle_K - \text{tr}(1_i^N u_{xx}(t_j, E_{t_j}, W_{E_{t_i}})Q) \Delta E_i \right) \right. \]

\[ \times \mathbb{E} \left( \mathbb{E} \left( \langle 1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) \delta W_j, \delta W_j \rangle_K \Big| G_{t_j} \right) - \text{tr}(1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q) \Delta E_j \right) \mathbb{E} \left( \mathbb{E} \left( \text{tr}(1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q) \Delta E_j \Big| G_{t_{j+1}} \right) \right) \}

\[ = \mathbb{E} \left\{ \left( \langle 1_i^N u_{xx}(t_j, E_{t_j}, W_{E_{t_i}}) \delta W_i, \delta W_i \rangle_K - \text{tr}(1_i^N u_{xx}(t_j, E_{t_j}, W_{E_{t_i}})Q) \Delta E_i \right) \right. \]

\[ \times \mathbb{E} \left( \mathbb{E} \left( \text{tr}(1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q) \Delta E_j \Big| G_{t_{j+1}} \right) \right) \}

\[ = 0. \]

Third, let \( f_{E_{t_{j+1}}, E_{t_j}}(\tau_1, \tau_2) \) be the joint density function of random variables, \( E_{t_{j+1}} \) and \( E_{t_j} \). Then, for \( t_{j+1} > t_j \), letting \( \mathbb{D} = \{(\tau_1, \tau_2) \in \)
\[ \mathbb{R}_+ \times \mathbb{R}_+ : 0 \leq \tau_2 \leq \tau_1 \}, \]

\[ \mathbb{E}\left( \langle 1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) \rangle W_j, \Delta W_j \rangle_K \right) \times \text{tr}(1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q) \Delta E_j \]

\[ = \int \int_D \mathbb{E}\left( \langle 1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) \rangle W_j, \Delta W_j \rangle_K \right) \times \text{tr}(1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q) \Delta E_j \]

\[ = \int \int_D \mathbb{E}\left( \text{tr}(1_j^N u_{xx}(t_j, \tau_2, W_{\tau_2})Q)^2 (\tau_1 - \tau_2)^2 \right) f_{E_{t_j+1}, E_{t_j}}(\tau_1, \tau_2) d(\tau_1, \tau_2) \]

\[ = \mathbb{E}\left( \text{tr}(1_j^N u_{xx}(t_j, E_{t_j}, W_{E_{t_j}})Q)^2 (\Delta E_j)^2 \right). \]

Finally, still for \( t_{j+1} > t_j \):

\[ \mathbb{E}\|\Delta W_j\|^4_K = \mathbb{E}\|W_{E_{t+1}} - W_{E_{t_j}}\|^4_K \]

\[ = \int \int_D \mathbb{E}\left( \|W_{\tau_1} - W_{\tau_2}\|^4_K \right) f_{E_{t+1}, E_{t_j}}(\tau_1, \tau_2) d(\tau_1, \tau_2) \]

\[ = 3(trQ)^2(\tau_1 - \tau_2)^2 f_{E_{t+1}, E_{t_j}}(\tau_1, \tau_2) d(\tau_1, \tau_2) \]

\[ = 3(trQ)^2\mathbb{E}(E_{t+1} - E_{t_j})^2 = 3(trQ)^2\mathbb{E}(\Delta E_j)^2. \]
Noting that \( u_{xx} \) is bounded on bounded subsets of \( H \), apply (3.18), (3.19), (3.20) and (3.21) to yield

\[
\mathbb{E}\left( \sum_{j=1}^{n-1} \left( \left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) (\Delta W_j), \Delta W_j \right)_K - tr\left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) Q \right) \Delta E_j \right)^2 \right)
\]

\[
= \sum_{j=1}^{n-1} \mathbb{E}\left( \left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) (\Delta W_j), \Delta W_j \right)_K - tr\left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) Q \right) \Delta E_j \right)^2 
\]

\[
+ \sum_{i \neq j=1}^{n-1} \mathbb{E}\left\{ \left( 1_N^j u_{xx}(t_i, E_{t_i}, W_{E_{t_i}}) (\Delta W_i), \Delta W_i \right)_K - tr\left( 1_N^j u_{xx}(t_i, E_{t_i}, W_{E_{t_i}}) Q \right) \Delta E_i \right\}
\times \left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) (\Delta W_j), \Delta W_j \right)_K - tr\left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) Q \right) \Delta E_j \right) \right\}
\]

\[
= \sum_{j=1}^{n-1} \left\{ \mathbb{E}\left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) (\Delta W_j), \Delta W_j \right)_K^2 - \mathbb{E}\left( tr\left( 1_N^j u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}) Q \right)^2 (\Delta E_j)^2 \right) \right\}
\]

\[
\leq \sup_{s \leq t, \|h\|_H \leq N} |u_{xx}(s, E_s, h)|^2_{\mathcal{L}(H)} \sum_{j=1}^{n-1} \left( \mathbb{E}\|\Delta W_j\|^4_K - (trQ)^2 \mathbb{E}(\Delta E_j)^2 \right)
\]

\[
= 2 \sup_{s \leq t, \|h\|_H \leq N} |u_{xx}(s, E_s, h)|^2_{\mathcal{L}(H)} (trQ)^2 \mathbb{E} \sum_{j=1}^{n-1} (\Delta E_j)^2 \to 0,
\]
where the convergence follows from the Fubini-Tonelli Theorem and the fact that the time-change $E_t$ has finite bounded variation. Additionally,

$$
P\left(\sum_{j=1}^{n-1} (1 - 1^N_j) \left\{ \langle u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}), \Delta W_j \rangle_K \right. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \right. \left. \r

Thus, it follows that

$$I_4 = \frac{1}{2} \sum_{j=1}^{n-1} \langle u_{xx}(t_j, E_{t_j}, W_{E_{t_j}}), (\Delta W_j, \Delta W_j)_K \rangle$$

$$\rightarrow \frac{1}{2} \int_0^t tr[u_{xx}(s, E_s, W_{E_s})Q]dE_s$$

in probability. Combining (3.15), (3.17) and (3.22), and taking the limit on the right hand side of (3.14) yields the Itô formula for the function $u(t, E_t, W_{E_t})$:

$$u(t, E_t, W_{E_t}) = u(0, 0, 0) + \int_0^t u_{t_1}(s, E_s, W_{E_s})ds + \int_0^t u_{t_2}(s, E_s, W_{E_s})dE_s$$

$$+ \frac{1}{2} \int_0^t tr[u_{xx}(s, E_s, W_{E_s})Q]dE_s$$

$$+ \int_0^t \langle u_x(s, E_s, W_{E_s}), dW_{E_s} \rangle_K,$$
in probability. Also note that
\begin{align*}
&u_{t_1}(t, E_t, k) = \langle F_x(X(0) + \psi t + \gamma E_t + \phi k), \psi \rangle_H \\
&u_{t_2}(t, E_t, k) = \langle F_x(X(0) + \psi t + \gamma E_t + \phi k), \gamma \rangle_H \\
&u_x(t, E_t, k) = \phi^* F_x(X(0) + \psi t + \gamma E_t + \phi k), \\
&u_{xx}(t, E_t, k) = \phi^* F_{xx}(X(0) + \psi t + \gamma E_t + \phi k),
\end{align*}
(3.24)
and

\begin{align*}
tr[F_{xx}(X(s))(\phi Q^{1/2})(\phi Q^{1/2})^*] &= tr[\phi^* F_{xx}(X(s))\phi Q].
\end{align*}
(3.25)

Substituting (3.24) and (3.25) into (3.23) yields
\begin{align*}
F(X(t)) &= F(X(0)) + \int_0^t \langle F_x(X(s)), \psi(s) \rangle_H ds \\
&+ \int_0^t \langle F_x(X(s)), \gamma(s) \rangle_H dE_s \\
&+ \int_0^t \langle F_x(X(s)), \phi(s) dW_E(s) \rangle_H \\
&+ \frac{1}{2} \int_0^t tr(F_{xx}(X(s))(\phi(s)Q^{1/2})(\phi(s)Q^{1/2})^*) dE_s
\end{align*}
(3.26)
in probability. Consequently, there is a subsequence such that the equality (3.26) holds almost surely. Therefore applying the second change of variable formula yields the desired result. \(\square\)

3.3. Stochastic differential equations (SDEs) driven by the time-changed \(Q\)-Wiener process

Let \(K\) and \(H\) be real separable Hilbert spaces, and let \(M(t) := W_{E_t}\) be a time-changed \(K\)-valued \(Q\)-Wiener process on a complete filtered probability space \((\Omega, \mathcal{G}, \{\mathcal{G}_{t \leq T}, \mathbb{P}\})\) with the filtration \(\mathcal{G}_t = \tilde{\mathcal{F}}_{E_t}\) satisfying the usual conditions. Consider the following type of
semilinear SDE driven by the time-changed $Q$-Wiener process on $[0,T]$ in $H$:

$$du(t) = Au(t)dt + B(u,t)dM(t)$$

with initial condition $u(0) = u_0$, where $A$ is the generator of a $C_0$-semigroup of operators $\{S(t), t \geq 0\}$.

**Theorem 3.9.** Assume that the following hypotheses are satisfied:

1) $S(t)$ is a contraction $C_0$-semigroup generated by $A$.

2) $B(t,u) : \mathcal{D}(\mathbb{R}^+, H) \to \tilde{\Lambda}_2(K_Q, H)$ is non-anticipating where $\mathcal{D}(\mathbb{R}^+, H)$ denotes the $H$-valued cadlag adapted processes with $\mathbb{R}^+$ as the time interval.

3) (Local Lipschitz property) For every $r > 0$, there exists a constant $K_r > 0$ such that for every $x,y \in \mathcal{D}(\mathbb{R}^+, H)$ and $t \geq 0$

$$\|B(t,x) - B(t,y)\|^2_{\mathbb{L}_2(K_Q,H)} \leq K_r \sup_{s \leq t} \|x(s) - y(s)\|^2_H$$

on the set $\{\omega : \sup_{s < t} \max(\|x(s,\omega)\|_H, \|y(s,\omega)\|_H) \leq r\}$.

4) (Growth condition) There exists a constant $K_0 > 0$ such that for every $x \in \mathcal{D}(\mathbb{R}^+, H)$ and $t \geq 0$,

$$\|B(t,x)\|^2_{\mathbb{L}_2(K_Q,H)} \leq K_0(1 + \sup_{s \leq t} \|x(s)\|^2_H).$$

5) $\mathbb{E}(\|u_0\|^2_H) < \infty$.

Then, the SDE (3.27) has a pathwise continuous solution of the form

$$u(t) = S(t)u_0 + \int_0^t S(t-s)B(s,u)dW_{E_s},$$

which is called the mild solution to the SDE.
Remark 3.10. Since $W_{E_t}$ is a Hilbert space-valued square integrable martingale, Theorem 3.9 follows from [31] where a similar result is provided for a SDE driven by a general Hilbert space-valued martingale.

Consider two types of SDEs in Hilbert space with the same initial condition $x_0$:

$$dY(t) = (AY(t) + F(t,Y(t)))dt + B(t,Y(t))dW_t,$$  \hspace{1em} (3.29)

and

$$dX(t) = (AX(t) + F(E_t,X(t)))dE_t + B(E_t,X(t))dW_{E_t}.$$  \hspace{1em} (3.30)

The following theorem establishes a deep connection between the classic SDE (3.29) and the time-changed SDE (3.30).

**Theorem 3.11.** (Duality of SDEs in Hilbert Space) Let $U_t$ be a $\beta$-stable subordinator and $E_t$ be the inverse of $U_t$, which is a finite $\tilde{F}_{E_t}$-measurable time-change.

- If an $H$-valued process $Y(t)$ satisfies SDE (3.29), then the $H$-valued process $X(t) := Y(E_t)$ satisfies SDE (3.30).

- If an $H$-valued process $X(t)$ satisfies SDE (3.30), then the $H$-valued process $Y(t) := X(U_{t^-})$ satisfies SDE (3.29).

**Proof.** First, consider the integral form of SDE (3.29):

$$Y(t) = x_0 + \int_0^t \left( AY(s) + F(s,Y(s)) \right) ds + \int_0^t B(s,Y(s))dW_s.$$  \hspace{1em} (3.31)
Suppose the $H$-valued process $Y(t)$ satisfies (3.31), and let $X(t) := Y(E_t)$. Applying the first change of variable formula gives

$$X(t) = x_0 + \int_0^t \left( AY(s) + F(s, Y(s)) \right) ds + \int_0^t B(s, Y(s)) dW_s$$

$$= x_0 + \int_0^t \left( AY_E(s) + F(E_s, Y(s)) \right) dE_s + \int_0^t B(E_s, Y(s)) dW_{E_s}$$

$$= x_0 + \int_0^t \left( AX(s) + F(E_s, X(s)) \right) dE_s + \int_0^t B(E_s, X(s)) dW_{E_s},$$

which is the corresponding integral form of SDE (3.30).

Similarly, suppose the $H$-valued process $X(t)$ satisfies the SDE (3.30). Using the second change of variable formula yields

$$X(t) = x_0 + \int_0^t \left( AX(s) + F(E_s, X(s)) \right) dE_s + \int_0^t B(E_s, X(s)) dW_{E_s}$$

$$= x_0 + \int_0^t \left( AX(U_{t-}) + F(E_{U_{t-}}, X(U_{t-})) \right) ds + \int_0^t B(E_{U_{t-}}, X(U_{t-})) dW_s$$

$$= x_0 + \int_0^t \left( AX(U_{t-}) + F(s, X(U_{t-})) \right) ds + \int_0^t B(s, X(U_{t-})) dW_s.$$

Let $Y(t) := X(U_{t-})$, then

$$Y(t) = x_0 + \int_0^t \left( AY(s) + F(s, Y(s)) \right) ds + \int_0^t B(s, Y(s)) dW_s,$$

which is the integral form of the SDE (3.29).

It is known from [5] that under appropriate conditions, the strong solution, $Y(t)$, of the SDE (3.29) exists and is unique in the form

$$Y(t) = x_0 + \int_0^t (AY(s) + F(s, Y(s))) ds + \int_0^t B(s, Y(s)) dW_s.$$
for all \( t \leq T \), \( \mathbb{P} \)-a.s. The existence and uniqueness of a strong solution to the time-changed SDE \( (3.30) \) is then established based on the duality theorem.

**Theorem 3.12.** Assume that the following hypotheses are satisfied:

1) \( W_t \) is a \( K \)-valued \( Q \)-Wiener process on a complete filtered probability space \((\Omega, \mathcal{F}, \{\mathcal{F}_t\}_{t \leq T}, \mathbb{P})\) with the filtration \( \mathcal{F}_t \) satisfying the usual conditions and \( E_t \) is the inverse \( \beta \)-stable subordinator which is independent of \( W_t \).

2) \( A \) is a linear bounded operator.

3) The coefficients \( F : \Omega \times [0,T] \times C([0,T], H) \to H \) and \( B : \Omega \times [0,T] \times C([0,T], H) \to \mathcal{L}_2(K_Q, H) \), where \( C([0,T], H) \) is the Banach space of \( H \)-valued continuous functions on \([0,T]\), satisfy the following conditions

   (a) \( F \) and \( B \) are jointly measurable, and for every \( 0 \leq t \leq T \), they are measurable with respect to the product \( \sigma \)-field \( \mathcal{F}_t \times \mathcal{C}_t \) on \( \Omega \times C([0,T], H) \), where \( \mathcal{C}_t \) is a \( \sigma \)-field generated by cylinders with bases over \([0,t]\).

   (b) There exists a constant \( L \) such that for all \( x \in C([0,T], H) \),

   \[
   \|F(\omega,t,x)\|_H + \|B(\omega,t,x)\|_{\mathcal{L}_2(K_Q,H)} \leq L(1 + \sup_{0 \leq s \leq T} \|x(s)\|_H)
   \]

   for \( \omega \in \Omega \) and \( 0 \leq t \leq T \).

   (c) For all \( x, y \in C([0,T], H) \), \( \omega \in \Omega \), \( 0 \leq t \leq T \), there exists \( K_0 > 0 \) such that

   \[
   \|F(\omega,t,x) - F(\omega,t,y)\|_H + \|B(\omega,t,x) - B(\omega,t,y)\|_{\mathcal{L}_2(K_Q,H)} \leq K_0 \sup_{0 \leq s \leq T} \|x(s) - y(s)\|_H.
   \]
4) \( \mathbb{E} \int_0^T \| B(t, Y(t)) \|^2_{L_2(K_0, H)} dt < \infty. \)

5) \( Y(t) \) is in the domain of \( A \text{d}P \times \text{d}t \)-almost everywhere.

6) \( x_0 \) is an \( \mathcal{F}_0 \)-measurable \( H \)-valued random variable.

Then, the time-changed SDE (3.30) has a unique strong solution, \( X(t) \), satisfying

\[
X(t) = x_0 + \int_0^t (AX(s) + F(E_s, X(s)))dE_s + \int_0^t B(E_s, X(s))dW_{Es}. \tag{3.32}
\]

**Proof.** From [5], based on conditions in Theorem 3.12, there is a unique solution \( Y(t) \) that satisfies the SDE (3.29). Moreover, it follows from Theorem 3.11 that \( X(t) := Y(E_t) \) satisfies the SDE (3.30). Therefore, there exists a solution to the time-changed SDE (3.30).

Now suppose there exists another solution to the SDE (3.30). Call this solution \( \hat{X}(t) \). Then, by Theorem 3.11, the process \( \hat{Y}(t) := \hat{X}(U_{t-}) \) is a solution to the SDE (3.29). Since the solution to the SDE (3.29) is unique from [5], it must be that \( \hat{Y}(t) = Y(t) \). Thus \( \hat{X}(U_{t-}) = Y(t) \) which implies that \( \hat{X}(t) = Y(E_t) = X(t) \). Therefore, the solution \( X(t) \) of the SDE (3.30) is unique and satisfies the desired integral equation (3.32). \( \square \)
4. Connections between Hilbert space-valued integrals
driven by time-changed $Q$-Wiener processes and
time-changed cylindrical Wiener processes,
respectively, and Walsh-type integrals

The objective of this section is to establish the equality of the three
integrals given below for appropriate integrands:

$$\int_0^T \int_{\mathbb{R}^N} g(t, x) M_E(dt, dx) = \int_0^T g(t)d\tilde{W}_{E_t} = \int_0^T \Phi_t^g \circ J^{-1} dW_{E_t},$$

where $M_E$ is a time-changed version of a worthy martingale measure,
$\tilde{W}_{E_t}$ is a time-changed version of a cylindrical Wiener process and
$W_{E_t}$ is a time-changed $Q$-Wiener process. Section 4.1 focuses on the
case of no time change. Section 4.2 provides the extension to the
time-changed case.

4.1. The case of no time change

The equality of these integrals in the case where there is no time
change are made in [7, 33]. The general idea is to first define a specific
random field $F$ and an associated Hilbert space $K$. After identifying
these objects, the integral with respect to the resulting cylindrical
Wiener process and the integral with respect to the developed mar-
tingale measure are shown to be the same for a particular class of
integrands. Also a connection between the cylindrical Wiener process
and a $Q$-Wiener process is made that leads to an equality of their
respective integrals. Thus, a combination of those results proves that
all three integrals are equal.
Let \( \{F(\phi) \mid \phi \in C^\infty_0(\mathbb{R}^+ \times \mathbb{R}^N)\} \) be a family of mean zero Gaussian random variables, called a Gaussian random field. The covariance of \( F \) is given by

\[
\mathbb{E}(F(\phi)F(\psi)) = \int_{\mathbb{R}^+} \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} \phi(s,x)f(x-y)\psi(s,y)dydxds, \tag{4.1}
\]

where \( f \) satisfies the following two conditions:

1) \( f \) is a non-negative, non-negative definite continuous function on \( \mathbb{R}^N\setminus\{0\} \), which is integrable in a neighborhood of 0;

2) for all \( \gamma \in \mathcal{S}(\mathbb{R}^N) \), the space of \( C^\infty \) functions which are rapidly decreasing along with all their derivatives, the following conditions hold:
   a) there exists a tempered measure \( \mu \) on \( \mathbb{R}^N \) such that for any \( m \in \mathbb{N}^+ \)
      \[
      \int_{\mathbb{R}^N} (1 + |\xi|^2)^{-m} \mu(d\xi) < \infty, \text{ and}
      \]
   b) \[
   \int_{\mathbb{R}^N} f(x)\gamma(x)dx = \int_{\mathbb{R}^N} F\gamma(\xi)\mu(d\xi), \text{ where } F\gamma \text{ is the Fourier transform of } \gamma.
   \]

For the remainder of Section 4, fix the specific random field \( F \) chosen in (4.1). We now define a Hilbert space associated with the fixed random field \( F \). Let \( K \) be the completion of the Schwartz space \( \mathcal{S}(\mathbb{R}^N) \) with the semi-inner product

\[
\langle \gamma, \psi \rangle_K := \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} \gamma(x)f(x-y)\psi(y)dydx
\]

\[
= \int_{\mathbb{R}^N} \mu(d\xi)F\gamma(\xi)\overline{F\psi(\xi)}, \tag{4.2}
\]
where $\gamma, \psi \in K$, and associated semi-norm $||\cdot||_K$. Then $K$ is a Hilbert space, see [36]. In the remaining of this section, the Hilbert space $K$ always refers to (4.2).

Moreover, after fixing a time interval $[0, T]$, it is possible to consider the set $K_T := L^2([0, T], ; K)$ with the norm

$$
\|g\|_{K_T}^2 = \int_0^T \|g(s)\|_K^2 \, ds.
$$

Note that $C_0^\infty([0, T] \times \mathbb{R}^N)$ is dense in $K_T$. It should also be noted that although $F$ was originally defined on smooth, compactly supported functions, it is possible to extend $F$ to functions of the form $1_{[0, T]}(\cdot) \varphi(\cdot)$ where $\varphi \in \mathcal{S}(\mathbb{R}^N)$. This follows since such an $F$ is a random linear functional such that $\gamma \mapsto F(\gamma)$ is an isometry from $(C_0^\infty([0, T] \times \mathbb{R}^N), \| \cdot \|_{K_T})$ into $L^2(\Omega, \mathcal{F}, \mathbb{P})$, i.e., a family of mean zero Gaussian random variables characterized by (4.1). For further details, see [33].

**Definition 4.1.** A cylindrical Wiener process on a Hilbert space $K$ as defined in (4.2) is a family of random variables $\{\tilde{W}_t, t \geq 0\}$ such that:

1. for each $h \in K$, $\{\tilde{W}_t(h), t \geq 0\}$ defines a Brownian motion with mean 0 and variance $t \langle h, h \rangle_K$;
2. for all $s, t \in \mathbb{R}^+$ and $h, g \in K$, $\mathbb{E}(\tilde{W}_s(h)\tilde{W}_t(g)) = (s \wedge t) \langle h, g \rangle_K$
   where $s \wedge t := \min(s, t)$.

From [33], the stochastic process $\{\tilde{W}_t, t \geq 0\}$ defined in terms of the fixed random field $F$ with covariance chosen in (4.1) is given
by
\[ \tilde{W}_t(\varphi) := F(1_{[0,t]}(\cdot)\varphi(\cdot)), \quad \text{for } \varphi \in K, \] (4.3)
is a cylindrical Wiener process on the Hilbert space \( K \). A complete orthonormal basis \( \{f_j\} \) can be chosen such that \( \{f_j\} \subset \mathcal{S}(\mathbb{R}^N) \) since \( \mathcal{S}(\mathbb{R}^N) \) is a dense subspace of \( K \). Consider the space \( L^2(\Omega \times [0,T];K) \) of predictable processes \( g \) such that
\[ \mathbb{E}\left( \int_0^T ||g(s)||_K^2 d\sigma \right) < \infty. \]

For \( g \in L^2(\Omega \times [0,T];K) \), the stochastic integral in Hilbert space \( H \) with respect to the cylindrical Wiener process \( \tilde{W}_t \) is defined as
\[ \int_0^T g(s)d\tilde{W}_s := \sum_{j=1}^{\infty} \int_0^T \langle g(s), f_j \rangle_K d\tilde{W}_s(f_j) \]
where particularly \( H = \mathbb{R} \) and \( \{f_j\} \) is an orthonormal basis of \( K \) in (4.2). The series is convergent in \( L^2(\Omega,\mathcal{F},\mathbb{P}) \), and the sum does not depend on the choice of orthonormal basis. Additionally, the following isometry holds:
\[ \mathbb{E}\left( \left( \int_0^T g(s)d\tilde{W}_s \right)^2 \right) = \mathbb{E}\left( \int_0^T ||g(s)||_K^2 d\sigma \right). \]

On the other hand, consider \( M_t \) defined by
\[ M_t(A) := F(1_{[0,t]}(\cdot)1_A(\cdot)), \quad t \in [0,T], A \in \mathcal{B}_b(\mathbb{R}^N), \] (4.4)
where \( F \) is the specific random field chosen in (4.1) and \( \mathcal{B}_b(\mathbb{R}^N) \) denotes the set of bounded Borel sets of \( \mathbb{R}^N \). The covariance of \( \{M_t(A)\} \)
is given by:

$$\mathbb{E}(M_t(A)M_t(B)) = \int_{\mathbb{R}^+} \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_{[0,t]}(s)1_A(x)f(x-y)1_B(y)d_zydxds$$

$$= t \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)d_zydx.$$

Therefore, $M_t$ is a martingale measure in the following sense.

**Definition 4.2.** A process $\{M_t(A)\}_{t \geq 0, A \in \mathcal{B}(\mathbb{R}^N)}$ is a martingale measure with respect to $\{\mathcal{F}_t\}_{t \geq 0}$ if:

1. for all $A \in \mathcal{B}(\mathbb{R}^N)$, $M_0(A) = 0$ a.s.;
2. for $t > 0$, $M_t$ is a sigma-finite $L^2(\mathbb{P})$-valued signed measure;
3. for all $A \in \mathcal{B}(\mathbb{R}^N)$, $\{M_t(A)\}_{t \geq 0}$ is a mean-zero martingale with respect to the filtration $\{\mathcal{F}_t\}_{t \geq 0}$.

Further, in order to define stochastic integrals with respect to a martingale measure, the martingale measure needs to be worthy.

**Definition 4.3.** A martingale measure $M$ is worthy if there exists a random sigma-finite measure $K(A \times B \times C, \omega)$, where $A, B \in \mathcal{B}(\mathbb{R}^N)$, $C \in \mathcal{B}(\mathbb{R}_+)$, and $\omega \in \Omega$, such that:

1. $A \times B \times C \mapsto K(A \times B \times C, \omega)$ is nonnegative definite and symmetric;
2. $\{K(A \times B \times (0,t])\}_{t \geq 0}$ is a predictable process for all $A, B \in \mathcal{B}(\mathbb{R}^N)$;
3. for all compact sets $A, B \in \mathcal{B}(\mathbb{R}^N)$ and $t > 0$,

$$\mathbb{E}(K(A \times B \times (0,t])) < \infty;$$
4. for all $A, B \in \mathcal{B}(\mathbb{R}^N)$ and $t > 0$,

$$|\mathbb{E}(M_t(A)M_t(B))| \leq \mathbb{K}(A \times B \times (0, t]) \text{ a.s.}$$

In particular, the stochastic integral with respect to a worthy martingale measure is defined in such a way that it is itself a martingale measure. First, consider elementary processes $g$ of the form

$$g(s, x, \omega) = 1_{[a, b]}(s)1_A(x)X(\omega) \quad (4.5)$$

where $0 \leq a < b \leq T$, $X$ is bounded and $\mathcal{F}_a$-measurable, and $A \in \mathcal{B}(\mathbb{R}^N)$. If $g$ is an elementary process as in (4.5), then define $g \cdot M$ by

$$g \cdot M_t(B) := \int_0^t \int_B g(s, x)M(ds, dx) = X(\omega) (M_{t\wedge b}(A \cap B) - M_{t\wedge a}(A \cap B)).$$

The definition of $g \cdot M$ can be extended by linearity to simple processes, which are finite sums of elementary processes. Let $\mathcal{P}_+$ denote the set of predictable processes $(\omega, t, x) \mapsto g(t, x; \omega)$ such that

$$\|g\|_+^2 := \mathbb{E} \left( \int_0^T \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} |g(t, x)| f(x - y) |g(t, y)| dy dx dt \right) < \infty.$$ 

Taking limits of simple processes, the definition of $g \cdot M$ extends to all $g \in \mathcal{P}_+$. From [7], $g \cdot M$ is a worthy martingale measure if $g \in \mathcal{P}_+$. Therefore, it makes sense to define the stochastic integral with respect to $M$ as a martingale measure in the following way:

$$\int_0^t \int_A g(s, x)M(ds, dx) =: g \cdot M_t(A).$$
The next result is Proposition 2.6 in [33]; the detailed proof is given there.

**Proposition 4.4.** ([33]) Suppose \( g \in \mathcal{P}_+ \). Then \( g \in L^2(\Omega \times [0,T];K) \) and

\[
\int_0^T \int_{\mathbb{R}^N} g(t,x)M(dt,dx) = \int_0^T g(t)d\tilde{W}_t,
\]

where \( M \) is the worthy martingale measure defined in (4.4) and \( \tilde{W}_t \) is the cylindrical Wiener process defined in (4.3).

Finally, Proposition 4.6 below provides conditions under which these integrals coincide with the integral with respect to a Q-Wiener process. Define the operator \( J : K \to K \) by

\[
J(h) := \sum_{j=1}^{\infty} \lambda_j^{1/2} \langle h, f_j \rangle_K f_j, \quad h \in K, \quad (4.6)
\]

where \( \{f_j\} \) is an orthonormal basis in \( K \) and \( \lambda_j \geq 0 \) satisfies

\[
\sum_{j=1}^{\infty} \lambda_j < \infty.
\]

**Lemma 4.5.** Let \( Q = JJ^* : K \to K \). Then \( Q \) has eigenvalues \( \lambda_j \) corresponding to the eigenvectors \( f_j \), i.e. \( Qf_j = \lambda_j f_j \).

**Proof.** Let \( f_j \) be the orthonormal basis as in (4.6), then

\[
Qf_j = (JJ^*)(f_j) = \sum_{l=1}^{\infty} \lambda_i^{1/2} \langle J^* f_j, f_i \rangle_K f_i = \sum_{l=1}^{\infty} \lambda_i^{1/2} \langle f_i, J f_j \rangle_K f_i
\]

\[
= \sum_{l=1}^{\infty} \lambda_i^{1/2} \langle f_j, \sum_{i=1}^{\infty} \lambda_i^{1/2} \langle f_i, f_i \rangle_K f_i \rangle_K f_i
\]

\[
= \sum_{i=1}^{\infty} \langle f_j, f_i \rangle_K f_i = \lambda_j f_j.
\]
Additionally, $Q$ is symmetric (self-adjoint), non-negative definite, and $trQ = \sum_{j=1}^{\infty} \lambda_j < \infty$. The operator $J : K \to K_Q$ is an isometry since

$$\|h\|_K = \|Q^{-1/2}J(h)\|_K = \|J(h)\|_{K_Q}, \ h \in K,$$

where $Q^{-1/2}$ denotes the pseudo-inverse of $Q^{1/2}$. Further, the inverse operator $J^{-1} : K_Q = Q^{1/2}K \to K$ is also an isometry. Therefore, taking $\{w_j(t)\}_{j=1}^{\infty}$ to be the family of independent Brownian motion processes defined by (4.3), a Q-Wiener process in $K$ as defined in Section 2 can be constructed as

$$W_t := \sum_{j=1}^{\infty} w_j(t)J(f_j) = \sum_{j=1}^{\infty} w_j(t)\lambda_j^{1/2}f_j. \quad (4.7)$$

As seen in Section 3 a predictable process $\phi$ will be integrable with respect to the Q-Wiener process $W_t$ if

$$\mathbb{E}\left( \int_0^T ||\phi(t)||_{L^2(K_Q,H)}^2 dt \right) < \infty, \quad (4.8)$$

Consider $g \in L^2(\Omega \times [0,T]; K)$, and define the operator, $\Phi^g_s : K \to \mathbb{R}$, by

$$\Phi^g_s(\eta) = \langle g(s), \eta \rangle_K, \ \eta \in K. \quad (4.9)$$

The following proposition is from Dalang and Quer-Sardanyons [33]; the proof is included here since it contains information that will be used later.
Proposition 4.6. If $g \in \mathcal{P}_+$ and $\Phi_t^g$ as defined in (4.9), then $\Phi_t^g \circ J^{-1}$ satisfies the condition (4.8) and

$$\int_0^T \Phi_t^g \circ J^{-1} dW_t = \int_0^T g(t) d\tilde{W}_t,$$

where $W_t$ is the $Q$-Wiener process defined in (4.7), $J$ is the operator defined in (4.6) and $\tilde{W}_t$ is the cylindrical Wiener process defined in (4.3).

Proof. First, to show that $\Phi_t^g \circ J^{-1} \in L_2(K_Q, \mathbb{R})$, note that

$$||\Phi_t^g \circ J^{-1}||_{L_2(K_Q, \mathbb{R})}^2 = \sum_{j=1}^{\infty} [(\Phi_t^g \circ J^{-1})(\lambda_j^{1/2} f_j)]^2 = \sum_{j=1}^{\infty} [\Phi_t^g(J^{-1} \lambda_j^{1/2} f_j)]^2$$

$$= \sum_{j=1}^{\infty} \langle g(t), J^{-1} Q^{1/2} f_j \rangle_k^2 = \sum_{j=1}^{\infty} \langle g(t), f_j \rangle_k^2$$

$$= ||g(t)||_K^2.$$ 

Therefore,

$$\mathbb{E} \left( \int_0^T ||\Phi_t^g \circ J^{-1}||_{L_2(K_Q, \mathbb{R})}^2 dt \right) = \mathbb{E} \left( \int_0^T ||g(t)||_K^2 dt \right) < \infty,$$

since $g \in \mathcal{P}_+$ implies that $g \in L^2(\Omega \times [0,T]; K)$. Thus, $\Phi_t^g \circ J^{-1}$ satisfies condition (4.8). Also, from Lemma 3.1

$$\int_0^T \Phi_t^g \circ J^{-1} dW_t = \sum_{j=1}^{\infty} \int_0^T (\Phi_t^g \circ J^{-1})(\lambda_j^{1/2} f_j) d\langle W_t, \lambda_j^{1/2} f_j \rangle_{K_Q}$$

$$= \sum_{j=1}^{\infty} \int_0^T \langle g(t), f_j \rangle_K d w_j(t)$$

$$= \sum_{j=1}^{\infty} \int_0^T \langle g(t), f_j \rangle_K d \tilde{W}_t(f_j)$$

$$= \int_0^T g(t) d\tilde{W}_t.$$
Therefore, combining Propositions 4.4 and 4.6 yields the desired integral connections.

**Corollary 4.7.** For $g \in \mathcal{P}_+$ and $\Phi^g_t$ as defined in (4.9),

$$
\int_0^T \int_{\mathbb{R}^N} g(t, x) M(dt, dx) = \int_0^T g(t) d\tilde{W}_t = \int_0^T \Phi^g_t \circ J^{-1} dW_t,
$$

where $M$ is the martingale measure defined in (4.4), $\tilde{W}_t$ is the cylindrical Wiener process defined in (4.3), $J$ is the operator defined in (4.6) and $W_t$ is the $Q$-Wiener process defined in (4.7).

4.2. The case of a time change

The previous section summarized results from [7, 33, 35] establishing the equality of integrals with respect to a martingale measure, a cylindrical Wiener process, and a $Q$-Wiener process. This section will extend those results to the time-changed case. The procedure for showing the equivalence of the integrals will be very similar to that used in the previous section. The same random field $F$ and Hilbert space $K$ are used to define time-changed versions of a cylindrical Wiener process and a martingale measure. Their associated integrals are then shown to be equal. Finally, a connection between the given time-changed cylindrical Wiener process and a time-changed $Q$-Wiener process leads to the equality of all three integrals.

First, recall that the time change, $E_t$, is the inverse of a $\beta$-stable subordinator. Then, the time-changed cylindrical Wiener process is defined as follows.
Definition 4.8. Let $K$ be a separable Hilbert space. A family of random variables $\{\tilde{W}_E_t, t \geq 0\}$ is a time-changed cylindrical Wiener process on $K$ if the following conditions hold:

1. for any $k \in K$, $\{\tilde{W}_E_t(k), t \geq 0\}$ defines a time-changed Brownian motion with mean 0 and variance $\mathbb{E}(E_t)(k,k)_K$; and
2. for all $s, t \in \mathbb{R}^+$ and $k, h \in K$,
   $$\mathbb{E}(\tilde{W}_E_s(k)\tilde{W}_E_t(h)) = \mathbb{E}(E_{s \wedge t})(k,h)_K.$$ 

Let $g : \mathbb{R}^+ \times \Omega \to K$ be any predictable process such that

$$\mathbb{E} \left( \int_0^T \|g(s)\|^2_K dE_s \right) < \infty. \quad (4.10)$$

Consider the series

$$\sum_{j=1}^{\infty} \int_0^T \langle g_s, f_j \rangle_K d\tilde{W}_E_s(f_j).$$

Convergence of this series in $L^2(\Omega, \mathcal{G}, \mathbb{P})$ is established in Proposition 4.9. This justifies defining the stochastic integral of $g$ with respect to a time-changed cylindrical Wiener process as follows:

$$\int_0^T g(s) d\tilde{W}_E_s := \sum_{j=1}^{\infty} \int_0^T \langle g_s, f_j \rangle_K d\tilde{W}_E_s(f_j).$$

Proposition 4.9. Let $g : \mathbb{R}^+ \times \Omega \to K$ in (4.2) be any predictable process satisfying condition (4.10). Then, the series

$$\sum_{j=1}^{\infty} \int_0^T \langle g_s, f_j \rangle_K d\tilde{W}_E_s(f_j) \quad (4.11)$$

converges in $L^2(\Omega, \mathcal{G}, \mathbb{P})$. 
Proof. Let

\[ Y_n := \sum_{j=1}^{n} \int_{T}^{0} \langle g_{s}, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j). \]

In order to show the convergence of the series defined in (4.11), it is sufficient to show \( \{Y_n\} \) is a Cauchy sequence in \( L^2(\Omega, \mathcal{G}, \mathbb{P}) \). For \( n > m \),

\[
||Y_n - Y_m||^2_2 = \mathbb{E}\left( \left( \sum_{j=m+1}^{n} \int_{T}^{0} \langle g_{s}, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j) \right)^2 \right)
\]

\[
= \mathbb{E}\left( \left( \sum_{j=m+1}^{n} \int_{T}^{0} \langle g_{s}, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j) \right) \left( \sum_{i=m+1}^{n} \int_{T}^{0} \langle g_{s}, f_i \rangle_{K} d\tilde{W}_{E_s}(f_i) \right) \right)
\]

\[
= \mathbb{E}\left( \sum_{j=m+1}^{n} \int_{T}^{0} \langle g_{s}, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j) \right)^2 \]

\[
+ \mathbb{E}\left( \sum_{j \neq i=m+1}^{n} \left[ \int_{T}^{0} \langle g_{s}, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j) \right] \left[ \int_{T}^{0} \langle g_{s}, f_i \rangle_{K} d\tilde{W}_{E_s}(f_i) \right] \right)
\]

\[ =: I + II. \]

Since the time-changed \( Q \)-Wiener process \( W_{E_t} \) defined in (2.4) is a square integrable martingale with respect to the filtration \( \mathcal{G}_t = \tilde{\mathcal{F}}_{E_t} \) defined in (2.5), for \( h = \lambda_j^{-1/2} f_j \in K, j = 1, 2, \cdots, \langle W_{E_t}, h \rangle_{K} \) is also a square integrable martingale with respect to the filtration \( \mathcal{G}_t \), i.e., for \( 0 < s < t \),

\[
\mathbb{E}(\langle W_{E_t}, h \rangle_{K}|\mathcal{G}_s) = \langle W_{E_s}, h \rangle_{k}.
\]

This implies that each projection, which is a time-changed Brownian motion, \( w_j(E_t), j = 1, 2, \cdots, \) is also a square integrable martingale with respect to the same filtration \( \mathcal{G}_t \), i.e., for \( 0 < s < t \),
$\mathbb{E}(w_j(E_t)|G_s) = w_j(E_s)$. Therefore, the integral,
\[
\int_0^T \langle g_s, f_j \rangle_K d w_j(E_s),
\]
is also a square integral martingale with respect to the filtration $G_t$, and it follows from the Itô isometry that
\[
\mathbb{E} \left[ \int_0^T \langle g_s, f_j \rangle_K d w_j(E_s) \right]^2 = \mathbb{E} \left[ \int_0^T \langle g_s, f_j \rangle_K^2 dE_s \right].
\]
Further, by a proof similar to that of Theorem 2.3, the product $w_i(E_t)w_j(E_t)$ is a square integrable martingale for $i \neq j$. This means the quadratic covariation process of martingales $w_i(E_t)$ and $w_j(E_t)$ is zero, i.e., $[w_i(E_t), w_j(E_t)] = 0$. Thus, using the martingale property of $w_j(E_s)$ and its associated integral, along with the Cauchy-Schwartz inequality
\[
I := \mathbb{E} \left( \sum_{j=i=m+1}^{n} \left[ \int_0^T \langle g_s, f_j \rangle_K d \tilde{W}_{E_s}(f_j) \right]^2 \right)
\]
\[
= \sum_{j=i=m+1}^{n} \mathbb{E} \left[ \int_0^T \langle g_s, f_j \rangle_K d w_j(E_s) \right]^2
\]
\[
= \sum_{j=i=m+1}^{n} \mathbb{E} \left[ \int_0^T \langle g_s, f_j \rangle_K^2 dE_s \right] \leq \sum_{j=i=m+1}^{\infty} \mathbb{E} \left[ \int_0^T \langle g_s, f_j \rangle_K^2 dE_s \right].
\]
By assumption (4.10),
\[
\sum_{j=1}^{\infty} \mathbb{E} \left[ \int_0^T \langle g_s, f_j \rangle_K^2 dE_s \right] = \mathbb{E} \left[ \int_0^T \| g_s \|_K^2 dE_s \right] < \infty,
\]
so the above tail of the partial sum in (4.12) converges to 0 as \( m \) (hence \( n \)) \( \to \infty \). Meanwhile,

\[
\begin{align*}
II &:= \mathbb{E}\left( \sum_{j \neq i = m+1}^{n} \left[ \int_{0}^{T} \langle gs, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j) \right] \left[ \int_{0}^{T} \langle gs, f_i \rangle_{K} d\tilde{W}_{E_s}(f_i) \right] \right) \\
&= \sum_{j \neq i = m+1}^{n} \mathbb{E}\left[ \left( \int_{0}^{T} \langle gs, f_j \rangle_{K} d\tilde{W}_{E_s}(f_j) \right) \left( \int_{0}^{T} \langle gs, f_i \rangle_{K} d\tilde{W}_{E_s}(f_i) \right) \right] \quad (4.13) \\
&= \sum_{j \neq i = m+1}^{n} \mathbb{E}\left( \int_{0}^{T} \langle gs, f_j \rangle_{K} \langle gs, f_i \rangle_{K} d[w_j(E_s), w_i(E_s)] \right) \\
&= 0.
\end{align*}
\]

Combining (4.12) and (4.13) yields \( ||Y_n - Y_m||_2^2 = I + II \to 0 \) as \( n, m \to \infty \). Thus, \( \{Y_n\} \) is a Cauchy sequence in \( L^2(\Omega, \mathcal{G}, \mathbb{P}) \), completing the proof.

The next proposition shows how to define a cylindrical process from the fixed random field \( F \) chosen in (4.1)

**Proposition 4.10.** For \( t \geq 0 \) and \( \phi \in K \), set

\[
\tilde{W}_{E_t}(\phi) := F(1_{[0,E_t]}(\cdot)\phi(*)). 
\quad (4.14)
\]

Then, the process \( \tilde{W}_{E_t} \) is a time-changed cylindrical Wiener process.

**Proof.** Consider a fixed \( \phi \in K \). \( \tilde{W}_{E_t}(\phi) \) is a time-changed Brownian motion by construction. Additionally,

\[
\begin{align*}
\mathbb{E}[\tilde{W}_{E_t}(\phi)] &= \mathbb{E}[F(1_{[0,E_t]}(\cdot)\phi(*))] = \int_{0}^{\infty} \mathbb{E}[F(1_{[0,\tau]}(\cdot)\phi(*))] f_{E_t}(\tau) d\tau \\
&= \int_{0}^{\infty} 0 \cdot f_{E_t}(\tau) d\tau = 0,
\end{align*}
\]
and
\[
\mathbb{E}[\tilde{W}_{E_t}(\phi)\tilde{W}_{E_t}(\phi)] = \mathbb{E}[\int_{[0,E_t]}(\phi(s))\int_{[0,E_t]}(\phi(s))]
\]
\[
= \mathbb{E}\left[\int_{\mathbb{R}^+} 1_{[0,E_t]}(s)\int_{\mathbb{R}^N} \phi(x)f(x-y)\phi(y)dydxds\right]
\]
\[
= \mathbb{E}\left[\langle \phi, \phi \rangle_K \int_{\mathbb{R}^+} 1_{[0,E_t]}(s)ds\right] = \mathbb{E}(E_t)\langle \phi, \phi \rangle_K.
\]

Further, for fixed \(s,t \in \mathbb{R}^+\) and \(\phi, \psi \in K\),
\[
\mathbb{E}[\tilde{W}_{E_t}(\phi)\tilde{W}_{E_s}(\psi)] = \mathbb{E}[\int_{[0,E_t]}(\phi(s))\int_{[0,E_s]}(\psi(s))]
\]
\[
= \mathbb{E}\int_{\mathbb{R}^+} 1_{[0,E_t]}(r)\int_{\mathbb{R}^N} \phi(x)f(x-y)\psi(y)dydxdydr
\]
\[
= \mathbb{E}\int_{\mathbb{R}^+} \langle \phi, \psi \rangle_K 1_{[0,E_t\wedge E_s]}(r)dr = \mathbb{E}(E_t\wedge s)\langle \phi, \psi \rangle_K.
\]

Therefore, according to the Definition 4.8, \(\tilde{W}_{E_t}\) is a time-changed cylindrical Wiener process.

Moreover, it follows from Proposition 4.9 that the integral with respect to the process \(\tilde{W}_{E_t}\) defined in (4.14) is well-defined.

On the other hand, \(E_t\) is independent of the martingale measure \(M_t(A)\). So, define a time-changed version of \(\{M_t(A)\}\) by
\[
M_{E_t}(A) = M(A \times [0, E_t]) := F(1_{[0,E_t]}(\cdot)1_A(*)).
\]
By conditioning on the time change, the covariance for \(M_{E_t}(A)\) is
\[
\mathbb{E}(M_{E_t}(A)M_{E_t}(B)) = \int_0^\infty \mathbb{E}(M_{E_t}(A)M_{E_t}(B))f_{E_t}(\tau)d\tau
\]
\[
= \int_0^\infty \tau \left(\int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx\right)f_{E_t}(\tau)d\tau
\]
\[
= \mathbb{E}(E_t) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx,
\]
where $f_{E_t}$ is the density function of $E_t$. Also

$$
\mathbb{E}[M_{E_t}(A)]^2 = \mathbb{E}(E_t) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_A(y)dydx < \infty
$$

for all $A \in \mathcal{B}_b(\mathbb{R}^N)$. Thus, $M_{E_t}(A)$ has a finite second moment for all $A \in \mathcal{B}_b(\mathbb{R}^N)$. Then the following theorem shows that \{\(M_{E_t}(A)\)\} is also a martingale measure.

**Theorem 4.11.** \{\(M_{E_t}(A)\)\}_{t \geq 0, A \in \mathcal{B}(\mathbb{R}^N)} is a martingale measure with respect to the filtration \{\(\tilde{F}_{E_t}\)\}_{t \geq 0}, where \(\tilde{F}_t\) in (2.5) is generated by the time change $E_t$ and independent Brownian motions $\tilde{W}_t(f_j), j = 1, 2, \cdots$ defined by (4.3).

**Proof.** It suffices to check the conditions in Definition (4.2). First, since $E_0 = 0$ a.s., $M_{E_0}(A) = M_0(A) = 0$ a.s. because $M_t(A)$ is a martingale measure.

Second, let $A, B \in \mathcal{B}(\mathbb{R}^N)$ be disjoint. Then, for fixed $\tau$, $M_\tau(A \cup B)$ and $M_\tau(A) + M_\tau(B)$ are mean zero Gaussian random variables and

$$
\text{Var}(M_\tau(A \cup B)) = \int \int f(x-y)dydx
$$

$$
= \int \int f(x-y)dydx + \int \int f(x-y)dydx + 2\int \int f(x-y)dydx
$$

$$
= \text{Var}(M_\tau(A)) + \text{Var}(M_\tau(B)) + 2\mathbb{E}(M_\tau(A)M_\tau(B))
$$

$$
= \text{Var}\left(M_\tau(A) + M_\tau(B)\right).
$$

Also note that

$$
M_\tau(A \cup B) = M_\tau(A) + M_\tau(B) \text{ a.s.}
$$
Thus, conditioning on the time change yields
\[
\mathbb{P}(M_{E_t}(A \cup B) = M_{E_t}(A) + M_{E_t}(B)) \\
= \int_0^\infty \mathbb{P}(M_\tau(A \cup B) = M_\tau(A) + M_\tau(B)) f_{E_t}(\tau) d\tau \\
= \int_0^\infty f_{E_t}(\tau) d\tau = 1,
\]
which means $M_{E_t}(\cdot)$ is additive a.s. Furthermore, assume $A_1 \supset A_2 \supset \ldots$ such that $\cap_n A_n = \emptyset$,
\[
\mathbb{E}[M_{E_t}(A_n)]^2 = \mathbb{E}(E_t) \int_{A_n \times A_n} f(x-y)dydx \rightarrow 0
\]
as $n \to \infty$, and so $M_{E_t}(A_n) \to 0$ in $L^2(\mathbb{P})$. This proves the countable additivity of $M_{E_t}(\cdot)$.

Finally, since $\{M_{E_t}(A)\}$ has a finite second moment, a similar argument to the proof of Theorem 2.6 shows that $\{M_{E_t}(A)\}$ is a martingale for all $A \in \mathcal{B}(\mathbb{R}^N)$. Thus, $\{M_{E_t}(A)\}$ is a martingale measure with respect to the filtration $\mathcal{G}_t = \tilde{\mathcal{F}}_{E_t}$. \(\square\)

Define the dominating measure $K$ by
\[
K(A \times B \times C) := \mathbb{E}(\lambda(\{E_s(\omega) : s \in C\})) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx,
\]for (4.16)
where $A, B \in \mathcal{B}(\mathbb{R}^N)$ and $\lambda$ is the Lebesgue measure on $C \in \mathcal{B}(\mathbb{R}_+)$. Then, the following theorem shows that the martingale measure $\{M_{E_t}(A)\}$ is worthy.
Theorem 4.12. The martingale measure \( \{ M_{E_t}(A) \}_{t \geq 0, A \in \mathcal{B}(\mathbb{R}^N)} \) is worthy with respect to the filtration \( \{ \tilde{\mathcal{F}}_{E_t} \}_{t \geq 0} \), i.e. the same one as defined in Theorem 4.11.

Proof. To show that the martingale measure \( \{ M_{E_t}(A) \} \) is worthy, it suffices to show that the dominating measure \( K \) defined in (4.16) satisfies the conditions of Definition 4.3.

1. For all \( C \in \mathcal{B}(\mathbb{R}^+) \),
\[
K(A \times B \times C) = \mathbb{E}(\lambda(\{ E_s(\omega) : s \in C \})) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx
\]
\[
= \mathbb{E}(\lambda(\{ E_s(\omega) : s \in C \})) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_B(x)f(x-y)1_A(y)dydx
\]
\[
= K(B \times A \times C).
\]
Additionally, for all \( A \in \mathcal{B}(\mathbb{R}^N), C \in \mathcal{B}(\mathbb{R}^+) \), since \( f \) is non-negative definite,
\[
K(A \times A \times C) = \mathbb{E}(\lambda(\{ E_s(\omega) : s \in C \})) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_A(y)dydx \geq 0.
\]

2. For all \( A, B \in \mathcal{B}(\mathbb{R}^N), t > 0, \)
\[
K(A \times B \times (0,t)) = \mathbb{E}(E_t) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx
\]
is \( \tilde{\mathcal{F}}_{E_t} \)-measurable.

3. For all compact sets \( A, B \in \mathcal{B}(\mathbb{R}^N) \) and \( t > 0, \)
\[
\mathbb{E}|K(A \times B \times (0,t))| = \mathbb{E}(E_t) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx < \infty.
\]
4. For all \( A, B \in \mathcal{B}(\mathbb{R}^N) \) and \( t > 0 \),

\[
|\mathbb{E}(M_{E_t}(A)M_{E_t}(B))| = \mathbb{E}(E_t) \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} 1_A(x)f(x-y)1_B(y)dydx
= K(A \times B \times (0,t]).
\]

Thus, the martingale measure \( \{M_{E_t}(A)\} \) is worthy. \( \square \)

As shown in the non-time-changed case, the stochastic integral with respect to a worthy time-changed martingale measure is also a worthy martingale measure. First, consider elementary processes \( g \) of the form

\[
g(s,x,\omega) = 1_{(a,b]}(s)1_A(x)X(\omega)
\quad (4.17)
\]

where \( 0 \leq a < b \leq T \), \( A \in \mathcal{B}(\mathbb{R}^N) \) and \( X \) is both bounded and \( \mathcal{G}_a := \tilde{\mathcal{F}}_{E_a} \)-measurable. For \( t, r \in [0,T] \), let

\[
M_{E_t \wedge E_r} := M_{E_{t \wedge r}},
\]

and let \( M_E(ds,dx) \) denote integration with respect to the martingale measure in the both \( s \) and \( x \). Then, define \( g \cdot M_E \) by

\[
g \cdot M_{E_t}(B) := X(\omega)(M_{E_t \wedge E_a}(A \cap B) - M_{E_t \wedge E_a}(A \cap B))
= \int_0^t \int_{\mathbb{R}^N} g(s,x)M_E(ds,dx).
\]

As usual, this definition of \( g \cdot M_E \) can be extended to finite sums of elementary processes and finally to predictable processes \( g \) such
that

\[ ||g||_1^2 := \mathbb{E}\left( \int_0^T \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} |g(t, x)| |f(x - y)| |g(t, y)| dy dx d\mathbb{E}_t \right) \]

\[ < \infty. \]

Let \( \mathcal{P}_1 \) denote the set of predictable processes \((\omega, t, x) \mapsto g(t, x; \omega)\) such that (4.18) holds. For \( g \in \mathcal{P}_1 \), \( g \cdot M_E \) is a worthy martingale measure and the stochastic integral with respect to \( M_E \) is defined by

\[ \int_0^t \int_A g(s, x) M_E(ds, dx) =: g \cdot M_E(t, A). \]

The following theorem connects an integral with respect to a time-changed martingale measure with an integral with respect to a time-changed cylindrical Wiener process.

**Theorem 4.13.** Let \( M_E \) be the time-changed martingale measure defined in (4.15) and \( \tilde{W}_E \) be the time-changed cylindrical process defined in (4.14). For \( g \in \mathcal{P}_1 \), then,

\[ \int_0^T \int_{\mathbb{R}^N} g(t, x) M_E(dt, dx) = \int_0^T g(t) d\tilde{W}_E(t). \]

**Proof.** First notice that since \( g \in \mathcal{P}_1 \),

\[ \mathbb{E}\left( \int_0^T ||g(s)||_1^2 d\mathbb{E}_s \right) = \mathbb{E}\left( \int_0^T \int_{\mathbb{R}^N} \int_{\mathbb{R}^N} g(s, x) f(x - y) g(s, y) dy dx d\mathbb{E}_s \right) \]

\[ \leq ||g||_1^2 < \infty, \]

which means \( g \) satisfies the condition (4.10). Further, since the set of elementary processes is dense in \( \mathcal{P}_1 \), it suffices to check that the
integrals coincide for elementary processes of the form

\[ g(s, x, \omega) = 1_{(a, b]}(s)1_A(x)X(\omega) \]

where \(0 \leq a < b \leq T\), \(A \in \mathcal{B}(\mathbb{R}^N)\) and \(X(\omega)\) is both bounded and \(\tilde{F}_{E_a}\)-measurable. For this, note that

\[
\begin{align*}
\int_0^T \int_{\mathbb{R}^N} g(t, x)M_E(dt, dx) &= X(M_{E_T \land E_b}(A) - M_{E_T \land E_a}(A)) \\
&= X(M_{E_b}(A) - M_{E_a}(A)) \\
&= X(F(1_{[0, E_b]}(\cdot)1_A(*)) - F(1_{[0, E_a]}(\cdot)1_A(*))) \\
&= X(F(1_{(E_a, E_b]}(\cdot)1_A(*))).
\end{align*}
\]

On the other hand, using the linearity of \(F\)

\[
\begin{align*}
\int_0^T g(t) d\tilde{W}_{E_t} &= \sum_{j=1}^{\infty} \int_a^b X(1_A, f_j)K d\tilde{W}_{E_t}(f_j) \\
&= X \sum_{j=1}^{\infty} \langle 1_A, f_j \rangle K(\tilde{W}_{E_b}(f_j) - \tilde{W}_{E_a}(f_j)) \\
&= X \sum_{j=1}^{\infty} \langle 1_A, f_j \rangle K[F(1_{[0, E_b]}(\cdot)f_j) - F(1_{[0, E_a]}(\cdot)f_j)] \\
&= X \sum_{j=1}^{\infty} \langle 1_A, f_j \rangle K[F(1_{(E_a, E_b]}(\cdot)f_j)] \\
&= X[F(1_{(E_a, E_b]}(\cdot) \sum_{j=1}^{\infty} \langle 1_A, f_j \rangle f_j] \\
&= X[F(1_{(E_a, E_b]}(\cdot)1_A(*))].
\end{align*}
\]

Next, a connection between the time-changed cylindrical Wiener process and the time-changed \(Q\)-Wiener process will be established.
Define

$$W_{E_t} := \sum_{j=1}^{\infty} w_j(E_t) J(f_j) \tag{4.19}$$

where $w_j(E_t) := \tilde{W}_{E_t}(f_j)$ are the time-changed Brownian motions defined in Proposition 4.5. Also for $g \in L^2(\Omega \times [0, T], K)$, define the operator $\Phi^g_\eta$ by

$$\Phi^g_\eta(\eta) := \langle g(s), \eta \rangle_K \eta \in K. \tag{4.20}$$

A predictable process $\phi$ will be integrable with respect to $W_{E_t}$ if

$$\mathbb{E} \left( \int_0^T \|\phi(t)\|_{L^2(K_Q, H)}^2 dE_t \right) < \infty. \tag{4.21}$$

The next result provides the connection between the integral with respect to the time-changed cylindrical Wiener process and the time-changed $Q$-Wiener process.

**Theorem 4.14.** Let $\tilde{W}_{E_t}$ be the time-changed cylindrical Wiener process as in (4.14) and let $W_{E_t}$ be the time-changed $Q$-Wiener process as in (4.19). Let $g \in \mathcal{P}_\dagger$. Then, $\Phi^g_s \circ J^{-1}$ satisfies condition (4.21) and

$$\int_0^T \Phi^g_s \circ J^{-1} dW_{E_s} = \int_0^T g(s) d\tilde{W}_{E_s}. \tag{4.21}$$

**Proof.** First, from the proof of Proposition 4.6, $\|\Phi^g_s \circ J^{-1}\|_{L^2(K_Q, K)}^2 = \|g(s)\|_{K}^2$. Thus,

$$\mathbb{E} \left( \int_0^T \|\Phi^g_s \circ J^{-1}\|_{L^2(K_Q, R)}^2 dE_t \right) = \mathbb{E} \left( \int_0^T \|g(t)\|_{K}^2 dE_t \right) < \infty,$$
since it was shown in Theorem 4.13 that $g \in \mathcal{P}_\dagger$ implies that condition (4.10) holds. So, $\Phi^g_t \circ J^{-1}$ satisfies the condition (4.21). Also from Definition 3.2

$$\int_0^T \Phi^g_s \circ J^{-1} dW_{Es} = \sum_{j=1}^\infty \int_0^T \Phi^g_s \circ J^{-1} (\lambda_j^{1/2} f_j) d\langle W_s, \lambda_j^{1/2} f_j \rangle_{KQ}$$

$$= \sum_{j=1}^\infty \int_0^T \langle g(s), f_j \rangle_{K} dw_j(E_s)$$

$$= \sum_{j=1}^\infty \int_0^T \langle g(s), f_j \rangle_{K} d\tilde{W}_{Es}(f_j)$$

$$= \int_0^T g(s) d\tilde{W}_{Es}.$$ 

Finally, combining the results of Theorems 4.13 and 4.14 yields the desired correspondence of integrals with respect to time-changed processes.

**Corollary 4.15.** For $g \in \mathcal{P}_\dagger$ and $\Phi^g_t$ as defined in (4.20),

$$\int_0^T \int_{\mathbb{R}^N} g(t, x) M_E(dt, dx) = \int_0^T g(t) d\tilde{W}_{Et} = \int_0^T \Phi^g_t \circ J^{-1} dW_{Et},$$

where $M_{Et}$ is the time-changed martingale measure as in (4.15), $\tilde{W}_{Et}$ is the time-changed cylindrical process as in (4.14), $J$ is the operator defined in (4.6) and $W_{Et}$ is the time-changed $Q$-Wiener process as in (4.19).
5. Fokker-Planck-Kolmogorov equations associated with the time-changed stochastic differential equations

In this section, the Fokker-Plank-Kolmogorov (FPK) equations corresponding to sub-diffusion processes in Hilbert space are introduced. A FPK equation is a deterministic differential equation whose solution is the probability density function for a stochastic process. A fundamental example of such an equation in finite dimensions is the heat equation, whose solution is the density for Brownian motion. These equations are important for several reasons. As in the heat equation, these equations are often helpful in understanding a scientific phenomenon. The connection of these equations to a stochastic process allows one to use information about the stochastic process to study these phenomena. Conversely, knowing FPK equations corresponding to a particular stochastic process is helpful in the simulation of this stochastic process, see [25, 26, 29]. More recently, the FPK equations corresponding to diffusion processes in infinite dimensional Hilbert spaces have been analyzed, see [15, 16, 17].

An introduction to the sub-diffusion processes on Hilbert spaces considered here and their corresponding FPK equations requires a preliminary discussion of diffusion processes on Hilbert space and their FPK equations. Consider the following classic SDE driven by the $Q$-Wiener process

\[
\begin{cases}
  dY(t) = [AY(t) + F(t, Y(t))]dt + CdW_t \\
  Y(0) = x \in H
\end{cases}
\]

(5.1)
where \( A : D(A) \subset H \to H \) is the infinitesimal generator of a \( C_0 \)-semigroup \( S(t) = e^{tA}, t \geq 0 \), in \( H \), and \( W_t \) is a \( K \)-valued \( Q \)-Wiener process on a complete filtered probability space \((\Omega, \mathcal{F}, \{\mathcal{F}_t\}_{t \leq T}, \mathbb{P})\) with the filtration \( \mathcal{F}_t \) satisfying the usual conditions. Suppose that \( F : \Omega \times [0, T] \times H \to H \), \( C : \Omega \times K \to H \) and \( C \in \Lambda_2(K_Q, H) \). Also assume the initial value \( x \) is an \( \mathcal{F}_0 \)-measurable, \( H \)-valued random variable. Let \( Y(t) \) be a strong solution of the SDE (5.1) so that \( Y(t) \) will satisfy the following integral equation:

\[
Y(t) = x + \int_0^t [AY(s) + F(s, Y(s))]ds + CW_t.
\]

The Kolmogorov operator \( L_0 \) corresponding to the classic SDE (5.1) is

\[
L_0 \phi(x) = \langle x, A^* D_x \phi(x) \rangle_H + \langle F(t, x), D_x \phi(x) \rangle_H + \frac{1}{2} \text{tr}[(CQ^{1/2})(CQ^{1/2})^* D_x^2 \phi(x)],
\]

where \( x \in H \), \( t \in [0, T] \), and \( D_x, D_x^2 \) denote the first- and second-order Fréchet derivatives in space, respectively. \( D(L_0) \) denotes the domain of the operator \( L_0 \) and \( A^* \) denotes the adjoint of the operator \( A \). More details on the domain \( D(L_0) \) are given in [15, 16, 17].

Let \( \mu(dt, dx) \) be a product measure on \([0, T] \times H\) of the type

\[
\mu(dt, dx) = \mu_t(dx)dt,
\]

where \( \mu_t \in \mathbb{P}(H) \) is a Borel probability measure on the Hilbert space \( H \) for all \( t \in [0, T] \). Let \( P_t^Y \) be the transition evolution operator on \( \mathcal{B}_b(H) \), the space of bounded, Borel-measurable functions on \( H \),
defined by

\[ P_t^Y \phi(x) = \mathbb{E}(\phi(Y(t))|Y(0) = x), \quad 0 \leq t \leq T, \quad \phi \in \mathcal{B}_b(H), \]  

and let \((P_t^Y)^*\) be its adjoint operator. Note that \((P_t^Y)\) is a semi-group generated by the Markov process \(Y(t)\). According to [17], it is possible to define the measure \(\mu_t^Y\) induced by the solution \(Y(t)\) as

\[ \mu_t^Y(dy) := (P_t^Y)^*\xi(dy), \tag{5.4} \]

where \(\xi \in \mathbb{P}(H)\) is the measure associated with the initial value \(x\).

The induced measure, \(\mu_t^Y(dy)\), is defined as

\[ \int_H \phi(y)\mu_t^Y(dy) = \int_H P_t^Y \phi(y)\xi(dy), \quad \text{for all } \phi \in \mathcal{B}_b(H). \tag{5.5} \]

Under the assumption

\[ \int_{[0,T] \times H} \left( |\langle y, A^* h \rangle_H| + \|F(t, y)\|_H \right) \mu(dt, dy) < \infty, \]

where \(h \in D(A^*)\), the induced measure \(\mu_t^Y\) satisfies the following FPK equation

\[ \frac{d}{dt} \int_H \phi(y)\mu_t^Y(dy) = \int_H L_0\phi(y)\mu_t^Y(dy), \quad \text{for } dt\text{-a.e., } t \in [0, T], \tag{5.6} \]

where the initial condition is

\[ \lim_{t \to 0} \int_H \phi(y)\mu_t^Y(dy) = \int_H \phi(y)\xi(dy). \tag{5.7} \]
Further, if the domain of the Kolmogorov operator $L_0$ is comprised of test functions, integration by parts yields

$$\frac{\partial}{\partial t} \mu_t^Y = L_0^* \mu_t^Y, \quad \mu_0^Y = \xi.$$  \hspace{1cm} (5.8)

Further details are given in [17].

The following lemma is needed to extend the FPK equations (5.6) and (5.7) or (5.8) associated to the solution of the classic SDE (5.1) to the case of the solution to an SDE driven by a time-changed $Q$-Wiener process.

**Lemma 5.1.** Let $U_\beta(t)$ be a $\beta$-stable subordinator with the cumulative distribution function $F_\tau(t) = \mathbb{P}(U_\beta(\tau) \leq t)$ and density function $f_\tau(t)$. Suppose the inverse of $U_\beta(t)$ is $E_t$ with the density function $f_{E_t}(\tau)$. Then, for any integrable function $h(\tau)$ on $(0, \infty)$, the function $q(t)$ defined by the following integral

$$q(t) := \int_0^{\infty} f_{E_t}(\tau) h(\tau) d\tau$$

has Laplace transform

$$\mathcal{L}_{t\rightarrow s}\{q(t)\} = s^{\beta-1}|\widehat{h(\tau)}|(s^\beta),$$

where $\widehat{h(\tau)}(s) = \mathcal{L}_{\tau\rightarrow s}\{h(\tau)\}$.

**Proof.** Using the self-similarity property of the $\beta$-stable subordinator $U_\beta(t)$, the distribution function $F_{E_t}(\tau)$ associated with the time
change $E_t$ is

$$F_{E_t}(\tau) = \mathbb{P}(E_t \leq \tau) = \mathbb{P}(U_\beta(\tau) > t) = 1 - \mathbb{P}(\tau^{1/\beta} U_\beta(1) \leq t) = 1 - \mathbb{P}\left(U_\beta(1) \leq \frac{t}{\tau^{1/\beta}}\right) = 1 - F_1\left(\frac{t}{\tau^{1/\beta}}\right). \tag{5.9}$$

Differentiating both sides of (5.9),

$$f_{E_t}(\tau) = -\frac{\partial}{\partial \tau}\left\{ \frac{1}{\tau^{1/\beta}} (Jf_1)\left(\frac{t}{\tau^{1/\beta}}\right) \right\},$$

where $J$ is an integral operator defined by

$$(Jf)(\frac{t}{a}) = \int_0^t f\left(\frac{s}{a}\right) ds \text{ for all } a > 0. \tag{5.10}$$

Therefore, the Laplace transform

$$\mathcal{L}_{t \to s}\left[\frac{1}{a}(Jf_1)\left(\frac{t}{a}\right)\right](s) = \frac{1}{a} \int_0^\infty (Jf_1)\left(\frac{t}{a}\right) e^{-st} dt = \frac{1}{as} \int_0^\infty f_1\left(\frac{t}{a}\right) e^{-st} dt = \frac{1}{s} \int_0^\infty f_1(\hat{t})e^{-as\hat{t}} d\hat{t} = \frac{1}{s} \tilde{f}_1(as), \tag{5.11}$$

where $\tilde{f}_1(s)$ is the Laplace transform of the function $f(t)$. Now consider

$$q(t) := \int_0^\infty f_{E_t}(\tau) h(\tau) d\tau = -\int_0^\infty \frac{\partial}{\partial \tau}\left\{ \frac{1}{\tau^{1/\beta}} (Jf_1)\left(\frac{t}{\tau^{1/\beta}}\right) \right\} h(\tau) d\tau.$$ 

Applying the Laplace transform of the $\beta$-stable subordinator, $U_\beta(t)$, in (2.2) and using (5.11) yields

$$\mathcal{L}_{t \to s}\{q(t)\} = -\int_0^\infty \frac{\partial}{\partial \tau}\left\{ \frac{1}{s} e^{-s\tau^\beta} \right\} h(\tau) d\tau = s^{\beta-1} [\widehat{h(\tau)}](s^\beta).$$

\[\square\]

Let $W_{E_t}$ be a timechanged $Q$-Wiener process on a complete filtered probability space $(\Omega, \mathcal{G}, \{\mathcal{G}_t\}_{t \leq T}, \mathbb{P})$ with the filtration $\mathcal{G}_t :=$
satisfying the usual conditions. Suppose that \( x \) is an \( \mathcal{F}_0 \)-measurable and \( H \)-valued random variable. Let the coefficients \( A, F, \) and \( C \) be the same as in the classic SDE (5.1). Consider the following autonomous SDE on \( H \) driven by \( W_{E_t} \) and on the time interval \([0, T]\):

\[
\begin{aligned}
\frac{dX(t)}{dt} &= (AX(t) + F(X(t)))dE_t + CdW_{E_t}, \\
X(0) &= x \in H, \ t \geq 0.
\end{aligned}
\]  

(5.12)

We derive the time-fractional FPK equation associated with the time-changed SDE (5.12) via two different methods: first by applying the time-changed Itô formula and second by using the duality in Theorem 3.11. The advantage of the first approach is that it directly reveals the connection between the time-fractional FPK equations in Theorems 5.2 and 5.3 and the time-changed SDE (5.12). The advantage of the second approach is that it reveals the connection between the time-fractional FPK equations in Theorems 5.2 and 5.3 and the classic FPK equations in (5.6) and (5.7) or (5.8).

**Theorem 5.2.** Suppose the coefficients \( A, F, \) and \( C \) of the time-changed SDE (5.12) satisfy the conditions in Theorem 3.12. Let \( X(t) \) be the solution to (5.12). Also suppose that \( X(U_\beta(t)) \) is independent of \( E_t \). Then the probability kernel \( \mu^X_t(dx) \) induced by the solution \( X(t) \) satisfies the following fractional integral equation

\[
D_0^\beta \int_H \phi(x) \mu^X_t(dx) = \int_H L_0 \phi(x) \mu^X_t(dx),
\]  

(5.13)
with initial condition $\mu_0^X(dx) = \xi(dx)$, where $\phi \in D(L_0)$, $L_0$ is the Kolmogorov operator defined in (5.2) and $D^\beta_t$ denotes the Caputo fractional derivative operator as defined in Theorem 2.4.

Proof. (Method via the time-changed Itô formula) Let $Y(t) := X(U_t)$. Since $E_t$ and $W_{E_t}$ are both constant on $[U_{t-},U_t]$, the integrals

$$\int_0^t (AX(s) + F(X(s)))dE_s \quad \text{and} \quad \int_0^t B(X(s))dW_{E_s}$$

are also constant on $[U_{t-},U_t]$. Therefore, since $X(t)$ is the solution to the SDE (5.12), $X(t)$ is also constant on $[U_{t-},U_t]$ and satisfies $Y(E_t) = X(U_{E_t}) = X(t)$. So, for a fixed time $t \in [0,T]$, let $\mu^X_t$ and $\mu^Y_t$ denote the probability measures induced on $H$ by the stochastic processes $X(t)$ and $Y(t)$, respectively. Thus, for $\phi \in D(L_0)$,

$$\mathbb{E}(\phi(X(t))) = \int_H \phi(x)\mu^X_t(dx). \quad (5.14)$$

Since $X(t) = Y(E_t)$, taking the expectation of $X(t)$ conditioned on $E_t$,

$$\mathbb{E}(\phi(X(t))) = \mathbb{E}(\phi(Y(E_t))) = \int_0^\infty \mathbb{E}(\phi(Y_\tau)|E_t = \tau)f_{E_t}(\tau)d\tau,$$

where $f_{E_t}(\tau)$ is the density function of $E_t$. By the assumption that $Y(t) = X(U_t)$ is independent of $E_t$,

$$\mathbb{E}(\phi(X(t))) = \int_0^\infty \int_H \phi(x)\mu^Y_\tau(dx)f_{E_t}(\tau)d\tau$$

$$= \int_H \phi(x)\int_0^\infty \mu^Y_\tau(dx)f_{E_t}(\tau)d\tau. \quad (5.15)$$
Since $\phi \in D(L_0)$ is arbitrary, combining (5.14) and (5.15) yields
\[
\mu^X_t(dx) = \int_0^\infty \mu^Y_\tau(dx) f_{E_t}(\tau) \, d\tau. \quad (5.16)
\]
Since $X(t)$ is constant on every interval $[U_{r-}, U_r]$, $X(U_{r-}) = X(U_r) = Y(r)$. Thus, by the time-changed Itô formula,
\[
\phi(X(t)) - \phi(x_0) = \int_0^{E_t} L_0 \phi(X(U_{r_-})) \, dr + \int_0^{E_t} \langle \phi_x(X(U_{r_-})) , CdW_r \rangle \\
= \int_0^{E_t} L_0 \phi(Y(r)) \, dr + \int_0^{E_t} \langle \phi_x(Y(r)) , CdW_r \rangle. \quad (5.17)
\]
Since $\phi \in D(L_0)$, the integral
\[
M(\tau) = \int_0^\tau \langle \phi_x(Y(r)) , CdW_r \rangle_H
\]
is a square integrable $\mathcal{F}_\tau$-martingale. Taking expectations on both sides of (5.17) and conditioning on $E_t$ gives
\[
\mathbb{E}[\phi(X(t))|X(0) = x_0] - \phi(x_0) \\
= \int_0^\infty \mathbb{E}[\int_0^\tau L_0 \phi(Y(r)) \, dr + M(\tau)|E_t = \tau, X(0) = x_0] f_{E_t}(\tau) \, d\tau \\
= \int_0^\infty \int_0^\tau \mathbb{E}[L_0 \phi(Y(r))|X(0) = x_0] \, dr f_{E_t}(\tau) \, d\tau \\
= \int_0^\infty \int_0^\tau \int_H L_0 \phi(x) \mu^Y_\tau(dx) \, dr f_{E_t}(\tau) \, d\tau \\
= \int_0^\infty \left( J P^Y(\tau) \right) f_{E_t}(\tau) \, d\tau, \quad (5.18)
\]
where $J$ is the integral operator as defined in (5.10) and $P^Y(r)$ is defined by
\[
P^Y(r) = \int_H L_0 \phi(x) \mu^Y_\tau(dx).
\]
On the other hand,

\[ E[\phi(X(t))|X(0) = x_0] - E[\phi(x_0)|X(0) = x_0] \]

\[ = \int_H \phi(x) \mu_t^X(dx) - \int_H \phi(x)\xi(dx), \tag{5.19} \]

which also implies the initial condition \( \mu_0^X(dx) = \xi(dx) \). Combining (5.18), (5.19), Lemma 5.1 yields

\[ L_t \to s \left\{ \int_H \phi(x) \mu_t^X(dx) \right\} - \frac{1}{s} \int_H \phi(x)\xi(dx) = s^{\beta - 1}[J_t \overline{PY}(\tau)](s^\beta) = \frac{s^{\beta - 1}}{s^\beta} [\overline{PY}(\tau)](s^\beta), \]

which, in turn, implies

\[ s^\beta L_t \to s \left\{ \int_H \phi(x) \mu_t^X(dx) \right\} - s^{\beta - 1} \int_H \phi(x)\xi(dx) = s^{\beta - 1}[\overline{PY}(\tau)](s^\beta), \tag{5.20} \]

Combining (5.16) and Fubini’s Theorem yields

\[ \int_H L_0 \phi(x) \mu_t^X(dx) = \int_H L_0 \phi(x) \int_0^\infty \mu_t^Y(dx) f_{E_t}(\tau)d\tau \]

\[ = \int_0^\infty \int_H L_0 \phi(x) \mu_t^Y(dx) f_{E_t}(\tau)d\tau \]

\[ = \int_0^\infty P_Y(\tau)f_{E_t}(\tau)d\tau. \tag{5.21} \]

Taking Laplace transforms of both sides of (5.21) gives

\[ L_t \to s \left\{ \int_H L_0 \phi(x) \mu_t^X(dx) \right\} = L_t \to s \left\{ \int_0^\infty P_Y(\tau)f_{E_t}(\tau)d\tau \right\} \]

\[ = s^{\beta - 1}[\overline{P_Y}(\tau)](s^\beta). \tag{5.22} \]

Recall that the following equality holds:

\[ L_t \to s \{ D_t^\beta f(t) \} = s^\beta L_t \to s \{ f(t) \} - s^{\beta - 1} f(0), \tag{5.23} \]
where \( f(t) \) is a real-valued function on \( t \geq 0 \). Therefore, combining (5.20) and (5.22) yields

\[
s^\beta L_{t \to s} \left\{ \int_H \phi(x) \mu_t^X (dx) \right\} - s^{\beta-1} \int_H \phi(x) \xi (dx) = L_{t \to s} \left\{ \int_H L_0 \phi(x) \mu_t^X (dx) \right\},
\]

which, together with (5.23), implies that

\[
D_t^\beta \int_H \phi(x) \mu_t^X (dx) = \int_H L_0 \phi(x) \mu_t^X (dx).
\]

\[\square\]

The next theorem gives the familiar differential form of the FPK equation for the solution to the time-changed SDE (5.12).

**Theorem 5.3.** Suppose the conditions in Theorem 5.2 hold. If the domain of the operator \( L_0 \) defined in (5.2) is a set of test functions, then the probability measure \( \mu_t^X \) induced by the solution \( X(t) \) satisfies the following time-fractional PDE

\[
D_t^\beta \mu_t^X = L_0^* \mu_t^X, \tag{5.24}
\]

with initial condition \( \mu_0^X (dx) = \xi (dx) \), where \( L_0^* \) is the adjoint of the operator \( L_0 \) and \( D_t^\beta \) denotes the Caputo fractional derivative operator.

**Proof.** The proof of Theorem 5.2 gives

\[
\int_H \phi(x) \mu_t^X (dx) - \int_H \phi(x) \xi (dx) = \int_0^\infty \int_0^\tau \int_H L_0 \phi(x) \mu_r^X (dx) dr f_E (\tau) d\tau.
\]
Since $\phi \in D(L_0)$ is a test function, applying the integration by parts operator yields

$$
\int_H \phi(x) \mu_t^X(dx) - \int_H \phi(x) \xi(dx) = \int_H \phi(x) \int_0^\infty \int_0^\tau L_0^* \mu_t^Y(dx) dr f_{E_t}(\tau) d\tau,
$$

which means

$$
\mu_t^X(dx) - \xi(dx) = \int_0^\infty \int_0^\tau L_0^* \mu_t^Y(dx) dr f_{E_t}(\tau) d\tau. \tag{5.25}
$$

Additional information on the adjoint operator, $L_0^*$, is given in [37].

The proof of Theorem 5.2 also gives

$$
\mu_t^X(dx) = \int_0^\infty \mu_t^Y(dx) f_{E_t}(\tau) d\tau. \tag{5.26}
$$

Taking the Laplace transforms in (5.25) and (5.26) gives

$$
s^\beta L_{t \to s} \left\{ \mu_t^X(dx) \right\} - s^{\beta-1} \xi(dx) = L_{t \to s} \left\{ L_0^* \mu_t^X(dx) \right\},
$$

which implies

$$
D_t^\beta \mu_t^X = L_0^* \mu_t^X,
$$

with initial condition $\mu_0^X(dx) = \xi(dx). \qed$

We now use the second approach based on duality to derive the FPK equation for the solution to the time-changed SDEs (5.12).

**Theorem 5.4.** Under the assumptions of Theorem 5.2, the time-fractional FPK equation associated with the time-changed SDE (5.12) follows from (5.13). Also if the domain of the operator $L_0$ defined in (5.2) is a set of test functions, then the time-fractional FPK equation has the form (5.24).
Proof. (Via duality) From the duality theorem, Theorem \[\text{3.11}\] the solution of the time-changed SDE \[\text{(5.12)}\] is

\[X(t) = Y(E_t),\]  \[(5.27)\]

where \(Y(t)\) is the solution to the classic SDE \[\text{(5.1)}\]. As in \[\text{(5.3)}\] and \[\text{(5.4)}\], define the transition evolution operator, \(P_t^X\), induced by the solution, \(X(t)\), as follows:

\[P_t^X \phi(x) = E(\phi(X(t))|X(0) = x_0) = E(\phi(Y(E_t))|X(0) = x_0)\]

\[= \int_0^\infty P_\tau^Y \phi(x) f_{E_t}(\tau) d\tau, \quad 0 \leq t \leq T, \quad \phi \in \mathcal{B}_b(H).\]  \[(5.28)\]

The probability measure, \(\mu_t^X(dx)\), induced by \(X(t)\) is

\[\mu_t^X(dx) := (P_t^X)^\ast \xi(dx),\]  \[(5.29)\]

which means for all \(\phi \in \mathcal{B}_b(H),\)

\[\int_H \phi(x) \mu_t^X(dx) := \int_H P_t^X \phi(x) \xi(dx) = \int_H \int_0^\infty P_\tau^Y \phi(x) \xi(dx) f_{E_t}(\tau) d\tau.\]  \[(5.30)\]

Therefore, the connection between the probability measures, \(\mu_t^X(dx)\) and \(\mu_t^Y(dy)\), is obtained from \[\text{(5.30)}\] by applying Fubini’s theorem

\[\int_H \phi(x) \mu_t^X(dx) = \int_0^\infty \int_H \phi(x) \mu_t^Y(dx) f_{E_t}(\tau) d\tau, \quad \text{for all } \phi \in \mathcal{B}_b(H).\]  \[(5.31)\]

i.e.,

\[\mu_t^X(dx) = \int_0^\infty \mu_t^Y(dx) f_{E_t}(\tau) d\tau, \quad \text{for all } \phi \in \mathcal{B}_b(H).\]  \[(5.32)\]
Appealing to Lemma 5.1 and taking Laplace transforms of both sides of (5.31) leads to

\[
\mathcal{L}_{t \to s} \left\{ \int_H \phi(x) \mu_t^X (dx) \right\} = s^{\beta-1} \mathcal{L}_{\tau \to s} \left\{ \int_H \phi(x) \mu_\tau^Y (dx) \right\} (s^\beta). \tag{5.33}
\]

On the other hand, taking Laplace transforms on both sides of (5.6) leads to

\[
s \mathcal{L}_{t \to s} \left\{ \int_H \phi(y) \mu_t^Y (dy) \right\} - \int_H \phi(y) \xi(dy) = \mathcal{L}_{t \to s} \left\{ \int_H L_0 \phi(y) \mu_t^Y (dy) \right\} \tag{5.34}
\]

Replacing \( s \) by \( s^\beta \) in (5.34) yields

\[
s^\beta \mathcal{L}_{t \to s} \left\{ \int_H \phi(y) \mu_t^Y (dy) \right\} (s^\beta) - \int_H \phi(y) \xi(dy)
= \mathcal{L}_{t \to s} \left\{ \int_H L_0 \phi(y) \mu_t^Y (dy) \right\} (s^\beta). \tag{5.35}
\]

Thus, combining (5.33) and (5.35) gives

\[
s^\beta \mathcal{L}_{t \to s} \left\{ \int_H \phi(x) \mu_t^X (dx) \right\} - s^{\beta-1} \int_H \phi(y) \xi(dy)
= s^{\beta-1} \mathcal{L}_{t \to s} \left\{ \int_H L_0 \phi(y) \mu_t^Y (dy) \right\} (s^\beta),
\]

which implies

\[
D_t^\beta \int_H \phi(x) \mu_t^X (dx) = \int_0^\infty \int_H L_0 \phi(y) \mu_\tau^Y (dy) f_{E_t} (\tau) d\tau
= \int_H L_0 \phi(y) \int_0^\infty \mu_\tau^Y (dy) f_{E_t} (\tau) d\tau \tag{5.36}
= \int_H L_0 \phi(y) \mu_t^X (dy),
\]

as required. Further, if the domain of the Kolmogorov operator, \( L_0 \), is comprised of test functions, taking Laplace transforms on both
sides of (5.38) yields

\[ \begin{align*}
    sL_{t \to s}\left\{ \mu^Y_t \right\} - \mu^Y_0 &= L_{t \to s}\left\{ L^*_0 \mu^Y_t \right\}, \quad (5.37)
\end{align*} \]

while taking Laplace transforms of both sides of (5.32) and applying Lemma 5.1 yields

\[ \begin{align*}
    L_{t \to s}\left\{ \mu^X_t \right\} &= s^{\beta-1}L_{t \to s}\left\{ \mu^Y_{\tau} \right\}(s^{\beta}). \quad (5.38)
\end{align*} \]

Finally, combining (5.37) and (5.38) yields

\[ D^\beta_t \mu^X_t = L^*_0 \mu^X_t. \]

\[ \square \]

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