## Chapter 2 <br> Gradient and divergence

### 2.1 Gradient

If our objective is to define a differential calculus on the Banach space $W$, why don't we use the notion of Fréchet derivative: A function $F: W \rightarrow \mathbf{R}$ is said to be Fréchet differentiable if there exists a continuous linear operator $A: W \rightarrow W$ such that

$$
\begin{equation*}
\lim _{\epsilon \rightarrow 0} \epsilon^{-1}\left\|F\left(\omega+\epsilon \omega^{\prime}\right)-F(\omega)-A \omega^{\prime}\right\|_{W}=0 \tag{2.1}
\end{equation*}
$$

for any $\omega \in W$ and any $\omega^{\prime} \in W$. In particular, a Fréchet differentiable function is continuous. One of the most immediate function we can think of is the so-called Itô map which sends a sample-path $\omega$ to the corresponding sample-path of the solution of a well defined stochastic differential equation. It is well known (see [Lej09, Section 3.3] for instance) that in dimension higher than one, this map is not continuous. This induces that the notion of Fréchet derivative is not well suited to a differential calculus on the Wiener space. Moreover, since we work on a probability space, measurable functions $F$ from $W$ into $\mathbf{R}$ are random variables, meaning that they are defined up to a negligible set. To avoid any inconsistency in a formula like (2.1), we must ensure that

$$
(F=G \mu \text { a.s. }) \Longrightarrow\left(F\left(.+\omega^{\prime}\right)=G\left(.+\omega^{\prime}\right) \mu \text { a.s. }\right)
$$

for any $\omega^{\prime}$. With the notations of Theorem 1.8, this requires that $T_{\omega^{\prime}}^{\#} \mu$ (the pushforward of the measure $\mu$ by the translation map $T_{\omega^{\prime}}$ ) to be absolutely continuous with respect to the Wiener measure $\mu$. This fact is granted only if $\omega^{\prime}$ belongs to $I_{1,2}$. These two reasons mean that we are to define the directional derivative of $F$ in a restricted class of possible perturbations.

Recall the diagram

and that $\mu$ is the Wiener measure on $\mathcal{W}$.
Definition 2.1. A function $F$ is said to be cylindrical if there exists an integer $n, f \in \operatorname{Schwartz}\left(\mathbf{R}^{n}\right)$, the Schwartz space on $\mathbf{R}^{n},\left(h_{1}, \cdots, h_{n}\right) \in \mathcal{H}^{n}$ such that

$$
F(\omega)=f\left(\delta h_{1}, \cdots, \delta h_{n}\right)
$$

The set of such functionals is denoted by $\mathcal{S}$.
Theorem 2.1. The set $\mathcal{S}$ is dense in $L^{p}(\mathrm{~W} \rightarrow \mathbf{R} ; \mu)$.
Proof. Let $\mathcal{D}_{n}$ be the dyadic subdivision of mesh $2^{-n}$ of $[0,1]$ and $\mathcal{F}_{n}=$ $\sigma\left\{B(t), t \in \mathcal{D}_{n}\right\}$. Since $B$ has continuous sample-paths, for any $t \in[0,1]$, there exists a sequence $\left(t_{n}, n \geq 0\right)$ such that $t_{n} \in \mathcal{D}_{n}$ for any $n$ and $t_{n} \rightarrow t$. Hence, $\vee_{n} \mathcal{F}_{n}=\mathcal{F}$ and the $L^{p}$ convergence theorem for martingales says that

$$
\mathbf{E}\left[F \mid \mathcal{F}_{n}\right] \xrightarrow[L^{p}]{n \rightarrow \infty} F .
$$

For $\epsilon>0$, let $n$ such that $\left\|F-\mathbf{E}\left[F \mid \mathcal{F}_{n}\right]\right\|_{L^{p}}<\epsilon$. The Doob Lemma entails that there exists $\phi_{n}$ measurable from $\mathbf{R}^{2^{n}}$ to $\mathbf{R}$ such that

$$
\mathbf{E}\left[F \mid \mathcal{F}_{n}\right]=\psi_{n}\left(B(t), t \in \mathcal{D}_{n}\right)
$$

where $t_{k}^{n}=k 2^{-n}$. Let $\mu_{n}$ be the distribution of the Gaussian vector $(B(t), t \in$ $\left.\mathcal{D}_{n}\right)$,

$$
\int\left|\psi_{n}\right|^{p} \mathrm{~d} \mu_{n}=\mathbf{E}\left[\left|\mathbf{E}\left[F \mid \mathcal{F}_{n}\right]\right|^{p}\right] \leq \mathbf{E}\left[|F|^{p}\right]<\infty
$$

That means that $\psi_{n}$ belongs to $L^{p}\left(\mathbf{R}^{n} \rightarrow \mathbf{R} ; \mu_{n}\right)$ hence for any $\epsilon>0$, there exists $\varphi_{\epsilon} \in \mathcal{S}\left(\mathbf{R}^{2^{n}}\right)$ such that $\left\|\psi_{n}-\varphi_{\epsilon}\right\|_{L^{p}\left(\mathbf{R}^{n} \rightarrow \mathbf{R} ; \mu_{n}\right)}<\epsilon$. Then, $\varphi_{\epsilon}(B(t), t \in$ $\mathcal{D}_{n}$ ) belongs to $\mathcal{S}$ and is within distance $2 \epsilon$ of $F$ in $L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$.

The gradient is first defined on cylindrical functionals.
Definition 2.2. Let $F \in \mathcal{S}, h \in \mathcal{H}$, with $F(\omega)=f\left(\delta h_{1}, \cdots, \delta h_{n}\right)$. Set

$$
\nabla F=\sum_{j=1}^{n} \partial_{j} F\left(\delta h_{1}, \cdots, \delta h_{n}\right) h_{j}
$$

so that

$$
\langle\nabla F, h\rangle_{\mathcal{H}}=\sum_{j=1}^{n} \partial_{j} F\left(\delta h_{1}, \cdots, \delta h_{n}\right)\left\langle h_{j}, h\right\rangle_{\mathcal{H}} .
$$

This definition is coherent with the natural definition of directional derivative.

## Lemma 2.1.

$$
\langle\nabla F, h\rangle_{\mathcal{H}}=\lim _{\epsilon \rightarrow 0} \frac{1}{\epsilon}(F(\omega+\epsilon h)-F(\omega))
$$

Proof. Due to the fact that $\delta h_{i}(\omega+h)=\delta h_{i}(\omega)+\left\langle h, h_{i}\right\rangle_{\mathcal{H}}$.
Example 2.1. For any $h \in \mathcal{H}$,

$$
\nabla_{h} B(t)=\langle h, t \wedge .\rangle=\int_{0}^{t} \dot{h}(s) \mathrm{d} s=\int_{0}^{1} \mathbf{1}_{[0, t]}(s) \dot{h}(s) \mathrm{d} s
$$

Remark that $\mathcal{S}$ is an algebra for the ordinary product.
Lemma 2.2. For $F \in \mathcal{S}, \phi \in \mathcal{C}^{\infty}$

$$
\begin{aligned}
\nabla_{h} F(\omega) & =\left.\frac{d}{d \epsilon} F(\omega+\epsilon h)\right|_{\epsilon=0} \\
\nabla(F G) & =F \nabla G+G \nabla F \\
\nabla \phi(F) & =\phi^{\prime}(F) \nabla F .
\end{aligned}
$$

Before going futher, we give/recall some elements about tensor products of Banach spaces. Let $X$ and $Y$ two Banach spaces, with respective dual $X^{*}$ and $Y^{*}$. For $x \in X$ and $y \in Y, x \otimes y$ is the bilinear form defined by:

$$
\begin{aligned}
x \otimes y: X^{*} \times Y^{*} & \longrightarrow \mathbf{R} \\
(\eta, \zeta) & \longmapsto\langle\eta, x\rangle_{X^{*}, X}\langle\zeta, y\rangle_{Y^{*}, Y}
\end{aligned}
$$

Definition 2.3 (See [Rya02, chapter 2]). The projective tensor product of $X$ and $Y$ is the completion of the vector space spanned by the finite linear combinations of some $x \otimes y$ for $x \in X$ and $y \in Y$, equipped with the norm

$$
\|z\|_{X \otimes Y}=\inf \left\{\sum_{i=1}^{n}\left\|x_{i}\right\|_{X}\left\|y_{i}\right\|_{Y}, z=\sum_{i=1}^{n} x_{i} \otimes y_{i}\right\}
$$

We need to take the infimum of all the possible representations of $z$ as a linear combinations of elementary tensor products since such a representation is by no means unique.

Example 2.2. One of the simplest situation we can imagine, is the tensor product of $L^{1}(\mathbf{R})$ by itself. The function

$$
\mathbf{1}_{[0,1]}(s) \otimes \mathbf{1}_{[0,2]}(t)+\mathbf{1}_{[1,2]}(s) \otimes \mathbf{1}_{[1,2]}(t)
$$

can be equally written as:

$$
\mathbf{1}_{[0,1]}(s) \otimes \mathbf{1}_{[0,1]}(t)+\mathbf{1}_{[0,2]}(s) \otimes \mathbf{1}_{[1,2]}(t)
$$

Proposition 2.1. For $X$ and $Y$ two reflexive Banach spaces, i.e. $\left(X^{*}\right)^{*}=X$. The dual of $W=X \otimes Y$ is the space $W^{*}=X^{*} \otimes Y^{*}$ with the duality pairing:

$$
\left\langle w^{*}, w\right\rangle_{W *, W}=\sum_{i, j}\left\langle x_{i}^{*}, x_{j}\right\rangle_{X *, X}\left\langle y_{i}^{*}, y_{j}\right\rangle_{Y *, Y}
$$

where $w=\sum_{j} x_{j} \otimes y_{j} \in X \otimes Y$ and $w^{*}=\sum_{i} x_{i}^{*} \otimes y_{i}^{*} \in X^{*} \otimes Y^{*}$. Moreover,

$$
\begin{align*}
& \left\|w^{*}\right\|_{W^{*}}=\sup _{\|w\|_{W}=1}\left|\left\langle w^{*}, w\right\rangle_{W^{*}, W}\right| \\
& =\sup \left\{\left|\left\langle w^{*}, x \otimes y\right\rangle_{W^{*}, W}\right|,\|x\|_{X}=1,\|y\|_{Y}=1\right\} . \tag{2.2}
\end{align*}
$$

Let $X$ be a Banach space and $\nu$ a measure on a space $E$. The set $L^{p}(E ; X, \nu)$ is the space of functions $\psi$ from $E$ into $X$ such that

$$
\int_{E}\|\psi(x)\|_{X}^{p} \mathrm{~d} \nu(x)<\infty
$$

Theorem 2.2 ([Rya02, page 30]). For $X$ a Banach space, the space $L^{p}(E \rightarrow \mathbf{R} ; \nu) \otimes X$ is isomorphic to $L^{p}(E \rightarrow X ; \nu)$.

Moreover, if $X=L^{p}(F \rightarrow \mathbf{R} ; \rho)$ then $L^{p}(E \rightarrow X ; \nu)$ is isometrically isomorphic to $L^{p}(E \times F \rightarrow \mathbf{R} ; \nu \otimes \rho)$. Moreover, the set of simple functions, i.e. functions of the form

$$
\sum_{j=1}^{n} f_{j}(s) \psi_{j}(x)
$$

where $f_{j} \in L^{p}(E \rightarrow \mathbf{R} ; \nu)$ and $\psi_{j} \in L^{p}(F \rightarrow \mathbf{R} ; \rho)$, is dense into $L^{p}(E \times F \rightarrow \mathbf{R} ; \nu \otimes \rho)$.

Theorem 2.3. For $F \in \mathcal{S}, \nabla F$ belongs to $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ for any $p \geq 1$.
Proof. Step 1. Assume $p>1$. Since

$$
L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu) \simeq L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu) \otimes \mathcal{H}
$$

we have

$$
\left(L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)\right)^{*} \simeq L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu) \otimes \mathcal{H}
$$

where $q=p /(p-1)$.
Step 2. Consider the set

$$
B_{q, \mathcal{H}}=\left\{(k, G) \in \mathcal{H} \times L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu),\|k\|_{\mathcal{H}}=1,\|G\|_{L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)}=1\right\} .
$$

Let $F=f(\delta h)$, for $p>1$, in view of Proposition 2.1, we have to prove that

$$
\sup _{(k, G) \in B_{q, \mathcal{H}}}\left|\langle\nabla F, k \otimes G\rangle_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu), L^{q}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}\right|<\infty
$$

By the very definition of the duality bracket,

$$
\begin{aligned}
\left|\langle\nabla F, k \otimes G\rangle_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu), L^{q}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}\right| & =\left|\mathbf{E}\left[\langle\nabla F, k\rangle_{\mathcal{H}} G\right]\right| \\
& =\left|\mathbf{E}\left[f^{\prime}(\delta h) G\right]\langle k, h\rangle_{\mathcal{H}}\right| \\
& \leq\left\|f^{\prime}\right\|_{\infty}\|G\|_{L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)}\|h\|_{\mathcal{H}}\|k\|_{\mathcal{H}} .
\end{aligned}
$$

Hence the supremum over $B_{q, \mathcal{H}}$ is finite. The same proof can be applied when $F=f\left(\delta h_{j}, 1 \leq j \leq m\right)$.
Step 2. For $p=1$, the previous considerations no longer prevail since an $L^{1}$ space is not reflexive so that we cannot apply (2.2). However, it is sufficient to see that $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ is included in $L^{1}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$.

It is an exercise left to the reader to see that the map

$$
\begin{aligned}
\operatorname{Id} \otimes I^{-1}: L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu) \otimes \mathcal{H} & \longrightarrow L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu) \otimes L^{2}([0,1] \rightarrow \mathbf{R} ; \lambda) \\
F \otimes h & \longmapsto \otimes \dot{h}
\end{aligned}
$$

is continuous. Moreover, Theorem 2.3 means for any $F \in \mathcal{S}, \nabla F$ belongs to $L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu) \otimes \mathcal{H}$. Hence there exists an element $\dot{\nabla} F$ of $L^{p}(\mathcal{W} \rightarrow$ $\mathbf{R} ; \mu) \otimes L^{2}([0,1] \rightarrow \mathbf{R} ; \lambda)$ such that

$$
\begin{aligned}
& \langle\nabla F, h\rangle_{\mathcal{H}}=\int_{0}^{1} \dot{\nabla}_{s} F \dot{h}(s) \mathrm{d} s \\
& \quad \text { and }\|F\|_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}=\mathbf{E}\left[\left(\int_{0}^{1}\left|\dot{\nabla}_{s} F\right|^{2} \mathrm{~d} s\right)^{p / 2}\right]^{1 / p} .
\end{aligned}
$$

Proposition 2.2 ([Yos95, page 77]). A map $T$ is closable if and only if

$$
\left(x_{n}, n \geq 0\right) \in \operatorname{Dom} T, \lim _{n} x_{n}=0 \text { and } \lim _{n} T x_{n}=y \Longrightarrow y=0
$$

Theorem 2.4. $\nabla$ is closable in $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ for $p>1$.
The proof is based on the following lemma which is crucial for the sequel.
Lemma 2.3 (Integration by parts). For $F$ and $G$ cylindrical, for $h \in \mathcal{H}$,

$$
\begin{equation*}
\mathbf{E}\left[G\langle\nabla F, h\rangle_{\mathcal{H}}\right]=-\mathbf{E}\left[F\langle\nabla G, h\rangle_{\mathcal{H}}\right]+\mathbf{E}[F G \delta h] . \tag{2.3}
\end{equation*}
$$

Proof. The Cameron-Martin theorem says that

$$
\mathbf{E}[F(\omega+\epsilon h) G(\omega+\epsilon h)]=\mathbf{E}\left[F(\omega) G(\omega) \exp \left(\delta h-\frac{1}{2}\|h\|_{\mathcal{H}}^{2}\right)\right]
$$

Differentiate both sides with respect to $\epsilon$, at $\epsilon=0$, to obtain

$$
\mathbf{E}\left[F\langle\nabla G, h\rangle_{\mathcal{H}}\right]+\mathbf{E}\left[G\langle\nabla F, h\rangle_{\mathcal{H}}\right]=\mathbf{E}[F G \delta h],
$$

which corresponds to Eqn. (2.3).
Proof (Proof of Theorem 2.4). If $F_{n}$ tends to 0 in $L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$ then the right-hand-side of Eqn. (2.3) tends to 0 . On the other hand, by definition of the convergence in $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$,

$$
\mathbf{E}\left[G\left\langle\nabla F_{n}, h\right\rangle_{\mathcal{H}}\right] \xrightarrow{n \rightarrow \infty}\langle\eta, h \otimes G\rangle_{L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu), L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)} .
$$

It means that for any $h \in \mathcal{H}$ and $G \in \mathcal{S}$,

$$
\begin{equation*}
\langle\eta, h \otimes G\rangle_{L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu), L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)}=0 . \tag{2.4}
\end{equation*}
$$

By density of $\mathcal{S}$ in $L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu),(2.4)$ holds for $G \in L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$. According to Theorem 2.2, $\langle\eta, \zeta\rangle_{L^{p}(W ; \mathcal{H}), L^{q}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}=0$ for any $\zeta \in L^{q}(\mathcal{W} \rightarrow$ $\mathcal{H} ; \mu)$, hence $\eta=0$.

Definition 2.4 (Gross-Sobolev derivative). A functional $F$ belongs to $\mathbb{D}_{p, 1}$ if there exists $\left(F_{n}, n \geq 0\right)$ which converges to $F \in L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$, such that $\left(\nabla F_{n}, n \geq 0\right)$ is Cauchy in $L^{p}(W ; \mathcal{H})$. Then, $\nabla F$ is defined as the limit of this sequence.

Remark 2.1. The closability of $\nabla$ ensures that the limit does not depend on the approximating sequence.

Lemma 2.4. Let $p>1$. Assume that there exists $\left(F_{n}, n \geq 0\right)$ which converges in $L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$ to $F$ such that $\sup _{n}\left\|\nabla F_{n}\right\|_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}$ is finite. Then, $F \in \mathbb{D}_{p, 1}$.

For this result, we need to invoke two theorems of functional analysis.
Definition 2.5 (Weak convergence). A sequence ( $x_{n}, n \geq 0$ ) is said to be weakly convergent in a Banach space $X$, if for every $\eta \in X^{*}$, $\left(\left\langle\eta, x_{n}\right\rangle_{X^{*}, X}, n \geq 0\right)$ is convergent.
Remark 2.2. Since $\left|\left\langle\eta, x_{n}-x\right\rangle_{X^{*}, X}\right| \leq\|\eta\|_{X^{*}}\left\|x_{n}-x\right\|_{X}$, strong convergence implies weak convergence but the converse is false. For instance, let ( $e_{n}, n \geq$ 0) a complete orthonormal basis in a Hilbert space $X$, on the one hand $\left\|e_{n}\right\|_{X}=1$. On the other hand, according to Parseval equality, for $\eta \in X^{*}=$ $X,\|\eta\|_{X}^{2}=\sum_{n}\left|\left\langle\eta, e_{n}\right\rangle_{X}\right|^{2}$. Hence, $\left(\left\langle\eta, x_{n}\right\rangle_{X^{*}, X}, n \geq 0\right)$ converges weakly to 0 . The convergence cannot hold in the strong sense.

Proposition 2.3 (Eberlein-Shmulyan,[Yos95, page 141]). Let $X$ be a reflexive Banach space, i.e. $\left(X^{*}\right)^{*}=X$. Then, any strongly bounded sequence admits a weakly convergent subsequence.

Remark 2.3. For any measure, $L^{p}$ spaces are reflexive only for $p \neq 1, \infty$. We do have that the dual of $L^{1}$ is $L^{\infty}$ but the dual of $L^{\infty}$ is larger than $L^{1}$.

Proposition 2.4 (Mazur, [Yos95, page 120]). Let $\left(x_{n}, n \geq 0\right)$ be a weakly convergent subsequence in a Banach space $X$ and set $x$ its limit. Then, for any $\epsilon>0$, there exist $n$ and $\left(\alpha_{i}, 1 \leq i \leq n\right)$ such that $\alpha_{i} \geq 0, \sum_{i} \alpha_{i}=1$ and

$$
\left\|\sum_{i=1}^{n} \alpha_{i} x_{n_{i}}-x\right\|_{X} \leq \epsilon
$$

Proof (Proof of Lemma 2.4). Since $\sup _{n}\left\|\nabla F_{n}\right\|_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}$ is finite, there exists a subsequence (see Proposition 2.3) which we still denote by $\left(\nabla F_{n}, n \geq\right.$ 0 ) weakly convergent in $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ to some limit denoted by $\eta$. For $k>0$, let $n_{k}$ be such that $\left\|F_{m}-F\right\|_{L^{p}}<1 / k$ for $m \geq n$. The Mazur's Theorem 2.4 implies that there exists a convex combination of elements of $\left(\nabla F_{m}, m \geq n_{k}\right)$ such that

$$
\left\|\sum_{i=1}^{M_{k}} \alpha_{i}^{k} \nabla F_{m_{i}}-\eta\right\|_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}<1 / k .
$$

Moreover, since the $\alpha_{i}^{k}$ are positive and sums to 1 ,

$$
\left\|\sum_{i=1}^{M_{k}} \alpha_{i}^{k} F_{m_{i}}-F\right\|_{L^{p}} \leq 1 / k
$$

We have thus constructed a sequence

$$
F^{k}=\sum_{i=1}^{M_{k}} \alpha_{i}^{k} F_{m_{i}}
$$

such that $F^{k}$ tends to $F$ in $L^{p}$ and $\nabla F^{k}$ converges in $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ to a limit. By the construction of $\mathbb{D}_{p, 1}$, this means that $F$ belongs to $\mathbb{D}_{p, 1}$ and that $\nabla F=\eta$.

Definition 2.6. The space $\mathbb{D}_{p, 1}$ is the closure of $\mathcal{S}$ for the norm

$$
\|F\|_{p, 1}=\mathbf{E}\left[|F|^{p}\right]^{1 / p}+\mathbf{E}\left[\|\nabla F\|_{\mathcal{H}}^{p}\right]^{1 / p}
$$

Corollary 2.1. Let $F$ belong to $\mathbb{D}_{p, 1}$ and $G$ to $\mathbb{D}_{q, 1}$ with $q=p /(p-1)$. If $h \in \mathcal{H}$, then Eqn. (2.3) holds:

$$
\mathbf{E}\left[G\langle\nabla F, h\rangle_{\mathcal{H}}\right]=-\mathbf{E}\left[F\langle\nabla G, h\rangle_{\mathcal{H}}\right]+\mathbf{E}[F G \delta h] .
$$

Proof. According to Lemma 2.3, it is true for $F$ and $G$ in $\mathcal{S}$. Let $\left(F_{n}, n \geq 0\right)$ a sequence of elements of $\mathcal{S}$ converging to $F$ in $\mathbb{D}_{p, 1}$. Since $G$ belongs to $\mathcal{S}$, $G$ and $\nabla_{h} G$ belong to $L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$. By Hölder inequality, we see that (2.3) holds for $F \in \mathbb{D}_{p, 1}$ and $G \in \mathcal{S}$. Repeat the same approach with an approximation of $G \in \mathbb{D}_{q, 1}$ by elements of $\mathcal{S}$.

We can now generalize the basic formulas to elements of $\mathbb{D}_{p, 1}$.
Theorem 2.5. For $F \in \mathbb{D}_{p, 1}$ and $G \in \mathbb{D}_{q, 1}$ (with $1 / p+1 / q=1 / r$ for $r>1$ ), for $\phi$ Lipschitz continuous, the product $F G$ belongs to $\mathbb{D}_{r, 1}$ and

$$
\begin{aligned}
& \nabla(F G)=F \nabla G+G \nabla F \\
& \nabla \phi(F)=\phi^{\prime}(F) \nabla F
\end{aligned}
$$

For the next theorem, we need to introduce the two families of projections: For any $t \in[0,1]$

$$
\begin{aligned}
\pi_{t}: \mathcal{H} & \longrightarrow \mathcal{H} & \dot{\pi}_{t}: L^{2} & \longrightarrow L^{2} \\
h=I^{1}(\dot{h}) & \longmapsto I^{1}\left(\dot{h} \mathbf{1}_{[0, t]}\right) & \dot{h} & \longmapsto \dot{h} \mathbf{1}_{[0, t]}
\end{aligned}
$$

We have

$$
\left\|\pi_{t} h\right\|_{\mathcal{H}}^{2}=\int_{0}^{1} \dot{h}(s)^{2} \mathbf{1}_{[0, t]}(s) \mathrm{d} s \leq\|\dot{h}\|_{L^{2}}^{2}=\|h\|_{\mathcal{H}}^{2}
$$

meaning that $\pi_{t}$ is a continuous contraction on $\mathcal{H}$. Moreover,

$$
\left.\left.\pi_{t}(s \wedge .)=I^{1}\left(\dot{\pi}_{t}\left(\mathbf{1}_{[0, s]}\right)\right)=I^{1}\left(\mathbf{1}_{[0, s]} \mathbf{1}_{[0, t]}\right)\right)=I^{1}\left(\mathbf{1}_{[0, s \wedge t]}\right)\right)=(t \wedge s) \wedge
$$

so that

$$
\pi_{t}(s \wedge .)= \begin{cases}s \wedge . & \text { if } s \leq t  \tag{2.5}\\ t \wedge . & \text { otherwise }\end{cases}
$$

Lemma 2.5. Let $F \in \mathbb{D}_{p, 1}$ and $\mathcal{F}_{t}=\sigma\{\omega(s)$, $s \leq t\}$. Then, $\mathbf{E}\left[F \mid \mathcal{F}_{t}\right]$ belongs to $\mathbb{D}_{p, 1}$ and we have

$$
\begin{equation*}
\pi_{t} \mathbf{E}\left[\nabla F \mid \mathcal{F}_{t}\right]=\nabla \mathbf{E}\left[F \mid \mathcal{F}_{t}\right] \tag{2.6}
\end{equation*}
$$

Furthermore, if $F$ is $\mathcal{F}_{t}$-measurable then $\dot{\nabla}_{s} F=0$ for all $s>t$.
Proof. Step 1. First consider that $F$ is cylindrical. For the sake of simplicity, imagine that

$$
F=f\left(B\left(t_{1}\right), B\left(t_{2}\right)\right) \text { with } t_{1}<t<t_{2}
$$

Then,

$$
\begin{align*}
\mathbf{E}\left[F \mid \mathcal{F}_{t}\right] & =\mathbf{E}\left[f\left(B\left(t_{1}\right), B\left(t_{2}\right)-B(t)+B(t)\right)\right] \\
& =\int_{\mathbf{R}} f\left(B\left(t_{1}\right), B(t)+x\right) p_{t-t_{2}}(x) \mathrm{d} x \\
& =\tilde{f}\left(B\left(t_{1}\right), B(t)\right), \tag{2.7}
\end{align*}
$$

where $p_{t-t_{2}}$ is the density of a centered Gaussian distribution of variance $\left(t_{2}-t\right)$ and

$$
\tilde{f}(u, v)=\int_{\mathbf{R}} f(u, v+x) p_{t-t_{2}}(x) \mathrm{d} x \text { belongs to } \operatorname{Schwartz}\left(\mathbf{R}^{2}\right) .
$$

On the one hand,

$$
\begin{equation*}
\nabla_{s} \mathbf{E}\left[F \mid \mathcal{F}_{t}\right]=\partial_{1} \tilde{f}\left(B\left(t_{1}\right), B(t)\right) t_{1} \wedge s+\partial_{2} \tilde{f}\left(B\left(t_{1}\right), B(t)\right) t \wedge s \tag{2.8}
\end{equation*}
$$

On the other hand,

$$
\begin{align*}
& \mathbf{E}\left[\nabla_{s} F \mid \mathcal{F}_{t}\right]=\mathbf{E}\left[\partial_{1} f\left(B\left(t_{1}\right), B\left(t_{2}\right)\right) \mid \mathcal{F}_{t}\right] t_{1} \wedge s \\
&+\mathbf{E}\left[\partial_{2} f\left(B\left(t_{1}\right), B\left(t_{2}\right)\right) \mid \mathcal{F}_{t}\right] t_{2} \wedge s . \tag{2.9}
\end{align*}
$$

The same reasoning as in (2.7) leads to

$$
\begin{align*}
\mathbf{E}\left[\partial_{i} f\left(B\left(t_{1}\right), B\left(t_{2}\right)\right) \mid \mathcal{F}_{t}\right]=\int_{\mathbf{R}} \partial_{i} f\left(B\left(t_{1}\right), B(t)\right. & +x) p_{t-t_{2}}(x) \mathrm{d} x \\
& =\partial_{i} \tilde{f}\left(B\left(t_{1}\right), B(t)\right) \tag{2.10}
\end{align*}
$$

for $i \in\{1,2\}$. In view of (2.10), Eqn. (2.9) becomes

$$
\begin{equation*}
\mathbf{E}\left[\nabla_{s} F \mid \mathcal{F}_{t}\right]=\sum_{i=1}^{2} \partial_{i} \tilde{f}\left(B\left(t_{1}\right), B(t)\right) t_{i} \wedge s \tag{2.11}
\end{equation*}
$$

Thus, according to (2.5),

$$
\begin{aligned}
\pi_{t} \mathbf{E}\left[\nabla_{s} F \mid \mathcal{F}_{t}\right] & =\sum_{i=1}^{2} \partial_{i} \tilde{f}\left(B\left(t_{1}\right), B(t)\right) \pi_{t}\left(t_{i} \wedge .\right)(s) \\
& =\sum_{i=1}^{2} \partial_{i} \tilde{f}\left(B\left(t_{1}\right), B(t)\right)\left(t_{i} \wedge t\right) \wedge s \\
& =\nabla_{s} \mathbf{E}\left[F \mid \mathcal{F}_{t}\right] .
\end{aligned}
$$

Step 2. For the general case, let $\left(F_{n}, n \geq 0\right)$ a sequence of elements of $\mathcal{S}$ converging to $F$ in $\mathbb{D}_{p, 1}$. We can construct a sequence of cylindrical functions which are $\mathcal{F}_{t}$ measurable and converge in $\mathbb{D}_{p, 1}$ to $\mathbf{E}\left[F \mid \mathcal{F}_{t}\right]$. For any $n$, there exist $t_{1}^{n}<\ldots<t_{k_{n}}^{n}$ such that $F_{n}=f_{n}\left(B\left(t_{1}^{n}\right), \cdots, B\left(t_{k_{n}}^{n}\right)\right)$. If $t_{j_{0}}^{n} \leq t<t_{j_{0}+1}^{n}$,
for $l \geq j_{0}+1$, replace $B\left(t_{l}^{n}\right)$ by

$$
\left(B\left(t_{l}^{n}\right)-B\left(t_{l-1}^{n}\right)\right)+\ldots+\left(B\left(t_{j_{0}+1}^{n}\right)-B(t)\right)+B(t) .
$$

Let $W^{n}$ the Gaussian vector whose coordinates are the independent Gaussian random variables $\left(B\left(t_{k_{n}}^{n}\right)-B\left(t_{k_{n}-1}^{n}\right), \cdots, B\left(t_{j_{0}+1}^{n}\right)-B(t)\right)$ and

$$
\begin{aligned}
\kappa_{n}: \mathbf{R}^{k_{n}} & \longrightarrow \mathbf{R}^{k_{n}} \\
w=\left(w_{i}, 1 \leq i \leq k_{n}\right) & \longmapsto w_{i} \text { if } i \leq j_{0}, \\
& \longmapsto w_{i}+B(t)+\sum_{l=1}^{i-j_{0}} W_{l}^{n} \text { if } i>j_{0} .
\end{aligned}
$$

Hence

$$
\mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]=\mathbf{E}\left[\left(f_{n} \circ \kappa_{n}\right)\left(B\left(t_{1}^{n}\right), \cdots, B\left(t_{j_{0}}\right)\right) \mid B\left(t_{1}^{n}\right), \cdots, B(t)\right]
$$

Starting from this identity, we can reproduce the latter reasoning and see that (2.6) holds for such functionals.
Step 3. It remains to prove that $\mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]$ converges to $F=\mathbf{E}\left[F \mid \mathcal{F}_{t}\right]$ in $\mathbb{D}_{p, 1}$. By Jensen inequality,

$$
\mathbf{E}\left[\left|\mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]-\mathbf{E}\left[F \mid \mathcal{F}_{t}\right]\right|^{p}\right] \leq \mathbf{E}\left[\left|F_{n}-F\right|^{p}\right] \xrightarrow{n \rightarrow \infty} 0 .
$$

According to Proposition 2.1, the dual of $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ is $L^{q}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ and

$$
\begin{aligned}
& \left\|\nabla \mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]-\nabla \mathbf{E}\left[F_{m} \mid \mathcal{F}_{t}\right]\right\|_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)} \\
= & \sup \left\{\left|\mathbf{E}\left[\left\langle\nabla \mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]-\nabla \mathbf{E}\left[F_{m} \mid \mathcal{F}_{t}\right], h \otimes G\right\rangle\right]\right|,\|h\|_{\mathcal{H}}=1,\|G\|_{L^{q}}=1\right\} .
\end{aligned}
$$

Then, (2.6) implies that

$$
\begin{aligned}
&\left|\mathbf{E}\left[\left\langle\nabla \mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]-\nabla \mathbf{E}\left[F_{m} \mid \mathcal{F}_{t}\right], h \otimes G\right\rangle\right]\right| \\
&=\left|\mathbf{E}\left[\left\langle\pi_{t} \nabla \mathbf{E}\left[F_{n}-F_{m} \mid \mathcal{F}_{t}\right], h\right\rangle_{\mathcal{H}} G\right]\right| \\
&=\left|\mathbf{E}\left[\left\langle\nabla\left(F_{n}-F_{m}\right), \pi_{t} h\right\rangle_{\mathcal{H}} \mathbf{E}\left[G \mid \mathcal{F}_{t}\right]\right]\right| \\
& \leq\left\|\nabla\left(F_{n}-F_{m}\right)\right\|_{L^{p}(W ; \mathcal{H})}\|h\|_{\mathcal{H}}\|G\|_{L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)} .
\end{aligned}
$$

Since $\left(\nabla F_{n}, n \geq 0\right)$ is a Cauchy sequence in $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$, so does the sequence ( $\nabla \mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right], n \geq 0$ ), hence it is a converging sequence. Since $\nabla$ is closable, the limit can only be $\nabla \mathbf{E}\left[F \mid \mathcal{F}_{t}\right]$.
Step 4. Let $H_{t}^{\perp}=\bigcap_{s \in[t, 1] \cap \mathbf{Q}} \operatorname{ker}\left(\epsilon_{s}-\epsilon_{t}\right)$; it is a denumerable intersection of closed subspaces of $\mathcal{H}$, hence it is closed in $\mathcal{H}$. By sample-paths continuity of the elements of $\mathcal{H}, \dot{h}(s)=0$ for $s>t$ means that $h(s)=h(t)$ for any $s>t$ and $s \in \mathbf{Q}$, which is equivalent to $h \in H_{t}^{\perp}$. There exists a subsequence, we
still denote by $\left(F_{n}, n \geq 0\right)$, such that $\nabla \mathbf{E}\left[F_{n} \mid \mathcal{F}_{t}\right]$ converges almost-surely in $\mathcal{H}$ to $\nabla \mathbf{E}\left[F \mid \mathcal{F}_{t}\right]$. Since $H_{t}^{\perp}$ is closed, $\nabla \mathbf{E}\left[F \mid \mathcal{F}_{t}\right]$ belongs to $H_{t}^{\perp}$.

As we saw above, an element $U$ of $L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ can be represented as

$$
\begin{equation*}
U(\omega, t)=\int_{0}^{t} \dot{U}(\omega, s) \mathrm{d} s, \text { for all } t \in[0,1] \tag{2.12}
\end{equation*}
$$

where $\dot{U}$ is measurable from $W \times[0,1]$ onto $\mathbf{R}$.
Definition 2.7. An $\mathcal{H}$-valued random variable is said to be adapted whenever there exists a process $\dot{U}$ adapted in the classical sense such that (2.12) holds.

We denote by $L_{a}^{2}(W ; \mathcal{H})$ the set of $\mathcal{H}$-valued adapted, random variables such that

$$
\mathbf{E}\left[\int_{0}^{1}|\dot{U}(s)|^{2} \mathrm{~d} s\right]=\mathbf{E}\left[\|U\|_{\mathcal{H}}^{2}\right]<\infty .
$$

It is a closed subspace of $L_{a}^{2}(W ; \mathcal{H})$.
Similarly $\mathbb{D}_{2,1}^{a}(\mathcal{H})$ is the subset of $L^{2}(W ; \mathcal{H})$ such that

$$
\mathbf{E}\left[\iint\left|\dot{\nabla}_{r} \dot{U}(s)\right|^{2} \mathrm{~d} r \mathrm{~d} s\right]=\mathbf{E}\left[\|\nabla U\|_{L^{2}([0,1] ; \mathcal{H})}^{2}\right]<\infty
$$

Proposition 2.5. Let $X$ and $Y$ two Banach spaces and $M$ a dense subset of $X$. Consider $\left(A_{n}, n \geq 0\right)$ a sequence of continuous linear maps from $X$ to $Y$ such that:

1. $\sup _{n \geq 0}\left\|A_{n}\right\|<\infty$,
2. for any $x \in M$, the sequence $\left(A_{n} x, n \geq 0\right)$ is Cauchy in $Y$.

Then, for any $x \in X$ and not only in $M$, the sequence $\left(A_{n} x, n \geq 0\right)$ is convergent in $F$ and the linear map defined by $A x=\lim _{n \rightarrow \infty} A_{n} x$ is continuous from $X$ into $Y$. Moreover, $\|A\| \leq \sup _{n}\left\|A_{n}\right\|$.

Theorem 2.6. Let $U$ belongs to $\mathbb{D}_{p, 1}^{a}(\mathcal{H})$, and $\mathcal{D}_{n}$ the dyadic partition of $[0,1]$ of step $2^{-n}$. Then,

$$
\dot{U}_{\mathcal{D}_{n}}(t)=\sum_{i=1}^{2^{n}-1} 2^{n}\left(\int_{(i-1) 2^{-n}}^{i 2^{-n}} \dot{U}(r) d r\right) \mathbf{1}_{\left(i 2^{-n},(i+1) 2-n\right]}(t)
$$

converges in $\mathbb{D}_{p, 1}^{a}(\mathcal{H})$ to $U$.
Proof. Since indicator functions with disjoint support are orthogonal in $L^{2}([0,1])$, we have

$$
\begin{aligned}
\int_{0}^{1}\left|\dot{U}_{\mathcal{D}_{n}}(t)\right|^{2} \mathrm{~d} t= & \sum_{i=1}^{2^{n}-1}\left(2^{n} \int_{(i-1) 2^{-n}}^{i 2^{-n}} \dot{U}(r) \mathrm{d} r\right)^{2} \int_{0}^{1} \mathbf{1}_{\left(i 2^{-n},(i+1) 2-n\right]}(t) \mathrm{d} t \\
& \leq \sum_{i=1}^{2^{n}-1} \int_{(i-1) 2^{-n}}^{i 2^{-n}}|\dot{U}(r)|^{2} \frac{\mathrm{~d} r}{2^{-n}} 2^{-n}=\int_{0}^{1}|\dot{U}(r)|^{2} \mathrm{~d} r
\end{aligned}
$$

according to the Jensen inequality. Hence,

$$
\mathbf{E}\left[\left(\int_{0}^{1}\left|\dot{U}_{\mathcal{D}_{n}}(t)\right|^{2} \mathrm{~d} t\right)^{p / 2}\right] \leq \mathbf{E}\left[\left(\int_{0}^{1}|\dot{U}(r)|^{2} \mathrm{~d} r\right)^{p / 2}\right]
$$

or in other words that the maps

$$
\begin{aligned}
\mathcal{D}_{n}: L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu) & \longrightarrow L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu) \\
U & I^{1}\left(\dot{U}_{\mathcal{D}_{n}}\right)
\end{aligned}
$$

are continuous and satisfy $\left\|\mathcal{D}_{n}\right\| \leq 1$. Let $M$ be the subset of $L^{p}(\mathcal{W} \rightarrow$ $\mathcal{H} ; \mu)$ composed of processes such that $\dot{U}$ has continuous sample-paths with $\mathbf{E}\left[\|\dot{U}\|_{\infty}^{p}\right]<\infty$. For such a process

$$
\begin{aligned}
\left\|\dot{U}-\dot{U}_{\mathcal{D}_{n}}\right\|_{L^{2}([0,1])}^{2} \leq & \sum_{i=1}^{2^{n}-1} \int_{(i-1) 2^{-n}}^{i 2^{-n}}\left(2^{n} \int_{(i-1) 2^{-n}}^{i 2^{-n}}|\dot{U}(r)-\dot{U}(t)| \mathrm{d} r\right)^{2} \mathrm{~d} t \\
& \leq \sum_{i=1}^{2^{n}-1} \int_{(i-1) 2^{-n}}^{i 2^{-n}} 2^{n} \int_{(i-1) 2^{-n}}^{i 2^{-n}}|\dot{U}(r)-\dot{U}(t)|^{2} \mathrm{~d} r \mathrm{~d} t
\end{aligned}
$$

by the Jensen inequality. Since $\dot{U}$ is a.s. continuous, for $t \in\left((i-1) 2^{-n}, i 2^{-n}\right.$,

$$
2^{n} \int_{(i-1) 2^{-n}}^{i 2^{-n}}|\dot{U}(r)-\dot{U}(t)|^{2} \mathrm{~d} r \xrightarrow[\text { a.s. }]{n \rightarrow \infty} 0
$$

Since $\mathbf{E}\left[\|\dot{U}\|_{\infty}^{p}\right]$ is finite, the dominated convergence theorem entails that

$$
\mathbf{E}\left[\left(\sum_{i=1}^{2^{n}-1} \int_{(i-1) 2^{-n}}^{i 2^{-n}} 2^{n} \int_{(i-1) 2^{-n}}^{i 2^{-n}}|\dot{U}(r)-\dot{U}(t)|^{2} \mathrm{~d} r \mathrm{~d} t\right)^{p / 2}\right] \xrightarrow{n \rightarrow \infty} 0
$$

Apply the Proposition 2.5 to see that $\dot{U}_{\mathcal{D}_{n}}$ tends to $\dot{U}$ in $L^{2}(\mathcal{W} \otimes[0,1] \rightarrow$ $\mathbf{R} ; \mu \otimes \lambda)$.

Remark that if $\dot{U}$ is adapted then so does $\dot{U}_{\mathcal{D}_{n}}$. Moreover, there exists a subsequence $\dot{U}_{\pi_{n_{k}}}$ which converges almost-surely to $\dot{U}$ and that guarantees the adaptability of $\dot{U}$.

If $U \in \mathbb{D}_{2,1}, \dot{\nabla}_{r} \dot{U}_{t}$ can be approximated in $L^{2}\left(\mathcal{W} \otimes[0,1]^{2} \rightarrow \mathbf{R} ; \mu \otimes \lambda^{\otimes 2}\right)$ by

$$
\sum_{i=1}^{2^{n}-1} 2^{n}\left(\int_{(i-1) 2^{-n}}^{i 2^{-n}} \dot{\nabla}_{r} \dot{U}(s) \mathrm{d} s\right) \mathbf{1}_{\left(i 2^{-n},(i+1) 2-n\right]}(t)
$$

Then, the same proof as before shows this approximation converges in $L^{2}\left(\mathcal{W} \otimes[0,1]^{2} \rightarrow \mathbf{R} ; \mu \otimes \lambda^{\otimes 2}\right)$ to $\dot{\nabla} \dot{U}$.

Theorem 2.7. For $U \in \mathbb{D}_{2,1}^{a}(\mathcal{H})$, the Itô integral of $\dot{U}$ belongs to $\mathbb{D}_{2,1}$ and for any $h \in \mathcal{H}$,

$$
\begin{align*}
\left\langle\nabla\left(\int \dot{U}(s) d B(s)\right)\right. & , h\rangle_{\mathcal{H}} \\
& =\int_{0}^{1} \dot{U}(s) \dot{h}(s) d s+\int_{0}^{1}\langle\nabla \dot{U}(s), h\rangle_{\mathcal{H}} d B(s) \tag{2.13}
\end{align*}
$$

converges in $\mathbb{D}_{2,1}^{a}(\mathcal{H})$ to $U$.
Proof. From the previous theorem, we know that $\langle\nabla \dot{U}(s), h\rangle_{\mathcal{H}}$ is adapted and square integrable so that its stochastic integral is well defined. For $U(t)=$ $U_{a} I^{1}\left(\mathbf{1}_{(a, b]}\right)(t)$ with $U_{a} \in \mathcal{F}_{a}$ and $U_{a} \in \mathbb{D}_{2,1}$, on the one hand, since $\nabla$ is a derivation operator, we have

$$
\begin{aligned}
\left\langle\nabla \left(\int \dot{U}(s) \mathrm{d}\right.\right. & B(s)), h\rangle_{\mathcal{H}} \\
& =\left\langle\nabla\left(U_{a}(B(b)-B(a))\right), h\right\rangle_{\mathcal{H}} \\
& =\left\langle\nabla U_{a}, h\right\rangle_{\mathcal{H}}(B(b)-B(a))+\int_{0}^{1} U_{a} \mathbf{1}_{(a, b]}(s) \dot{h}(s) \mathrm{d} s \\
& =\int_{0}^{1}\left\langle\nabla U_{a}, h\right\rangle_{\mathcal{H}} \mathbf{1}_{(a, b]}(s) \mathrm{d} B(s)+\int_{0}^{1} U_{a} \mathbf{1}_{(a, b]}(s) \dot{h}(s) \mathrm{d} s \\
& =\int_{0}^{1} \dot{U}(s) \dot{h}(s) \mathrm{d} s+\int_{0}^{1}\langle\nabla \dot{U}(s), h\rangle_{\mathcal{H}} \mathrm{d} B(s) .
\end{aligned}
$$

By linearity, Eqn. (2.13) holds for simple processes as in Theorem 2.6. Since for $U$ with continuous sample-paths, $U_{\pi}$ tends in $L^{2}(W \times[0,1], \mu \otimes \lambda)$ to $U$, in virtue of Lemma 2.4, it remains to prove that

$$
\sup _{\pi} \mathbf{E}\left[\left\|\nabla \int \dot{U}_{\pi}(s) \mathrm{d} B(s)\right\|_{\mathcal{H}}^{p}\right]<\infty
$$

Consider the map

$$
\begin{aligned}
\int_{0}^{1} \nabla \dot{U}(s) \mathrm{d} B(s): \mathcal{H} & \longrightarrow L_{a}^{2}(\mu) \\
& h \longmapsto\left\langle\int_{0}^{1} \nabla \dot{U}(s), h\right\rangle_{\mathcal{H}} \mathrm{d} B(s) .
\end{aligned}
$$

We expect to compute

$$
\mathbf{E}\left[\sup _{\|h\|_{\mathcal{H}}=1}\left|\int_{0}^{1}\langle\nabla \dot{U}(s), h\rangle_{\mathcal{H}} \mathrm{d} B(s)\right|^{p}\right] .
$$

Without the supremum inside the expectation, we would usually refer to the Burkholder-Davis-Gundy inequality. In order to deal with the supremum, we note that

$$
t \longmapsto \int_{0}^{t} \nabla \dot{U}(s) \mathrm{d} B(s)
$$

is an Hilbert valued martingale and it satisfies also a BDG inequality:

$$
\begin{aligned}
& \mathbf{E}\left[\sup _{t \leq 1}\left|\int_{0}^{t} \nabla \dot{U}(s) \mathrm{d} B(s)\right|^{p}\right] \leq c_{p} \mathbf{E}\left[\left(\int_{0}^{1}\|\nabla \dot{U}(s)\|_{\mathcal{H}}^{2} \mathrm{~d} s\right)^{p / 2}\right] \\
& \quad=c_{p} \mathbf{E}\left[\left(\int_{0}^{1} \int_{0}^{1}\left|\dot{\nabla}_{r} \dot{U}(s)\right|^{2} \mathrm{~d} r \mathrm{~d} s\right)^{p / 2}\right]=c_{p}\|\nabla U\|_{L^{p}(W ; \mathcal{H} \otimes \mathcal{H})}^{p}
\end{aligned}
$$

Combining (2.13) with this upper-bound, we get

$$
\mathbf{E}\left[\left\|\nabla \int \dot{U}_{\pi}(s) \mathrm{d} B(s)\right\|_{\mathcal{H}}^{p}\right] \leq c\left(\left\|U_{\pi}\right\|_{L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}^{p}+\left\|\nabla U_{\pi}\right\|_{L^{p}(W ; \mathcal{H} \otimes \mathcal{H})}^{p}\right)
$$

We conclude with Theorem 2.6.
For cylindrical functions, we can clearly define higher order derivative following the same rule. The only difficulty is to realize that the second (respectively $k$-th) order gradient belongs to $\mathcal{H}^{\otimes(2)}$ (respectively $\left.\mathcal{H}^{\otimes(k)}\right)$ : For instance, for $F=f\left(\delta h_{j}, 1 \leq j \leq n\right)$,

$$
\begin{aligned}
&\left\langle\nabla^{(2)} F, h \otimes k\right\rangle=\sum_{j, l=1}^{n} \partial_{j, l} f\left(\delta h_{j}, 1 \leq j \leq n\right)\left\langle h_{j}, h\right\rangle_{\mathcal{H}} \otimes\left\langle h_{l}, k\right\rangle_{\mathcal{H}} \\
&=\left\langle\nabla\left(\langle\nabla F, h\rangle_{\mathcal{H}}\right), k\right\rangle_{\mathcal{H}}
\end{aligned}
$$

Definition 2.8. For any $p>1$ and $k \geq 1, \mathbb{D}_{p, k}$ is the closure of $\mathcal{S}$ with respect to the norm

$$
\|F\|_{p, k}=\|F\|_{p}+\sum_{j=1}^{k}\left\|\nabla^{(j)} F\right\|_{L^{p}\left(W ; \mathcal{H}^{\otimes(j)}\right)}
$$

The space of test functions is $\mathbb{D}=\cap_{p>1} \cap_{k \geq 1} \mathbb{D}_{p, k}$.

### 2.2 Divergence

For a matrix $M \in \mathcal{M}_{n, p}(\mathbf{R})$, its adjoint, which turns to coincide with its transpose, is defined by the identity:

$$
\langle M x, y\rangle_{\mathbf{R}^{p}}=\left\langle x, M^{*} y\right\rangle_{\mathbf{R}^{n}}
$$

We see that to define an adjoint, we need to have a notion a scalar product or more generally of a duality bracket. This means that if $M$ is continuous from a Banach $E$ into a Banach $F$, its adjoint is a continuous map from $F^{*}$ into $E^{*}$ defined by the identity:

$$
\langle M x, y\rangle_{F, F^{*}}=\left\langle x, M^{*} y\right\rangle_{E, E^{*}}
$$

For any $q>1$, the Gross-Sobolev derivative, which we denoted by $\nabla$, is continuous between the two spaces:

$$
\mathbb{D}_{q, 1} \subset L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu) \longrightarrow L^{q}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)
$$

Therefore its adjoint is a map from

$$
\begin{aligned}
\left(L^{q}(W \rightarrow \mathcal{H} ; \mu)\right)^{*}=L^{p}(\mathcal{W} & \rightarrow \mathcal{H} ; \mu) \\
& \longrightarrow\left(L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)\right)^{*}=L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)
\end{aligned}
$$

with $1 / p+1 / q=1$ and must satisfy the identity

$$
\begin{aligned}
\langle\nabla F, U\rangle_{L^{q}(\mathcal{W} \rightarrow \mathcal{H} ; \mu), L^{p}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)} & =\left\langle F, \nabla^{*} U\right\rangle_{L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu), L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)} \\
\Longleftrightarrow \mathbf{E}\left[\langle\nabla F, U\rangle_{\mathcal{H}}\right] & =\mathbf{E}[F \delta U]
\end{aligned}
$$

An additional difficulty comes from the fact that $\nabla$ is not defined on the whole of $L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$ but only on the subset $\mathbb{D}_{q, 1}$, hence we need to take some restrictions in the definition of the adjoint.
Definition 2.9. Let $p>1$. Let $\operatorname{Dom}_{p} \nabla^{*}$ be the set of $\mathcal{H}$-valued random variables $U$ for which there exists $c_{p}(U)$ such that for any $F \in \mathbb{D}_{q, 1}$,

$$
\left|\mathbf{E}\left[\langle\nabla F, U\rangle_{\mathcal{H}}\right]\right| \leq c_{p}(U)\|F\|_{L^{q}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)}
$$

In this case, we define $\nabla^{*} U$ as the unique element of $L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$ such that

$$
\mathbf{E}\left[\langle\nabla F, U\rangle_{\mathcal{H}}\right]=\mathbf{E}\left[F \nabla^{*} U\right] .
$$

Remark 2.4 ( $\nabla^{*}$ coincides with the Wiener integral on $\mathcal{H}$ ). Recall that $\delta$ is the Wiener integral. We now show that $\delta=\left.\nabla^{*}\right|_{\mathcal{H}}$. For any $F \in \mathcal{S}$, according to (2.3), we have

$$
\begin{equation*}
\mathbf{E}\left[\langle\nabla F, h\rangle_{\mathcal{H}}\right]=\mathbf{E}[F \delta h] \tag{2.14}
\end{equation*}
$$

and $\delta h$ is a Gaussian random variable of variance $\|h\|_{\mathcal{H}}^{2}$, thus belongs to any $L^{q}(\mathcal{H} \rightarrow \mathbf{R} ; \mu)$ for any $q>1$. Hence,

$$
\left|\mathbf{E}\left[\langle\nabla F, h\rangle_{\mathcal{H}}\right]\right| \leq c\|h\|_{\mathcal{H}}\|F\|_{L^{p}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)} .
$$

This means that $h$ belongs to $\operatorname{Dom}_{p} \delta$ and (2.14) entails that $\nabla^{*} h=\delta h$. Henceforth, in the following, we will use the notation $\delta$ instead of $\nabla^{*}$ and we keep for further reference the fundamental formula

$$
\begin{equation*}
\mathbf{E}\left[\langle\nabla F, U\rangle_{\mathcal{H}}\right]=\mathbf{E}[F \delta U] \tag{2.15}
\end{equation*}
$$

for any $F \in \mathbb{D}_{q, 1}$ and $U \in \operatorname{Dom}_{p} \delta$.
In usual deterministic calculus, if $a$ is a constant, then trivially

$$
\begin{equation*}
\int a u(s) \mathrm{d} s=a \int u(s) \mathrm{d} s \tag{2.16}
\end{equation*}
$$

For Itô integrals, this property does not hold any longer since we may have a problem of adaptability: If $a$ is a random variable, not belonging to $\mathcal{F}_{0}$ and $u$ is an adapted process with all the required integrability properties, then the process $(a u(s), s \geq 0)$ is not adapted so that $\int a u(s) \mathrm{d} B(s)$ is not even well defined. For the divergence, since we got rid of the adaptability hypothesis, we can prove a formula analog to (2.16) which is a simple consequence of the fact that $\nabla$ is a derivation operator.

Theorem 2.8. Let $U \in \operatorname{Dom}_{p} \delta$ and $a \in \mathbb{D}_{q, 1}$ with $1 / p+1 / q=1 / r$. Then, $a U \in \operatorname{Dom}_{r} \delta$ and

$$
\begin{equation*}
\delta(a U)=a \delta U-\langle\nabla a, U\rangle_{\mathcal{H}} \tag{2.17}
\end{equation*}
$$

Proof. Step 1. We first prove that the right-hand-side belongs to $L^{r}(\mathcal{W} \rightarrow$ $\mathbf{R} ; \mu)$.

$$
\begin{equation*}
\mathbf{E}\left[|a \delta U|^{r}\right] \leq \mathbf{E}\left[|a|^{p}\right]^{r / p} \mathbf{E}\left[|\delta U|^{q}\right]^{r / q} \tag{2.18}
\end{equation*}
$$

and

$$
\begin{align*}
\mathbf{E}\left[\left|\langle\nabla a, U\rangle_{\mathcal{H}}\right|^{r}\right] & \leq \mathbf{E}\left[\|\nabla a\|_{\mathcal{H}}^{r}\|U\|_{\mathcal{H}}^{r}\right] \\
& \leq \mathbf{E}\left[\|\nabla a\|_{\mathcal{H}}^{q}\right]^{r / q} \mathbf{E}\left[\mid U \|_{\mathcal{H}}^{p}\right]^{r / p} \\
& \leq\|a\|_{\mathbb{D}_{p, 1}}^{r}\|U\|_{\mathbb{D}_{q, 1}}^{r} . \tag{2.19}
\end{align*}
$$

Step 2. Denote $r^{*}=r /(r-1)$. For $F \in \mathbb{D}_{r^{*}, 1}$, since $\nabla$ is a true derivation,

$$
\begin{align*}
\mathbf{E}\left[\langle\nabla F, a U\rangle_{\mathcal{H}}\right] & =\mathbf{E}\left[\langle a \nabla F, U\rangle_{\mathcal{H}}\right] \\
& =\mathbf{E}\left[\langle\nabla(a F)-F \nabla a, U\rangle_{\mathcal{H}}\right]  \tag{2.20}\\
& =\mathbf{E}[F a \delta U]-\mathbf{E}\left[F\langle\nabla a, U\rangle_{\mathcal{H}}\right]
\end{align*}
$$

According to (2.18) and (2.19), (2.20) implies that

$$
\left|\mathbf{E}\left[\langle\nabla F, a U\rangle_{\mathcal{H}}\right]\right| \leq\|a\|_{\mathbb{D}_{p, 1}}\|U\|_{\mathbb{D}_{q, 1}}\|F\|_{L^{r^{*}}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)}
$$

Hence, $a U$ belongs to $\mathrm{Dom}_{r} \delta$.
Step 3. At last, (2.20) implies (2.17) by identification.
We have already seen that the Itô integral coincides with the Wiener integral for deterministic integrands provided that we identify $h$ and $\dot{h}$. We now show that modulo the same identification, the divergence of adapted processes coincides with their Itô integral.

Corollary 2.2 (Divergence extends Itô integral). Let $U \in \mathbb{D}_{2,1}^{a}(\mathcal{H})$. Then, $U$ belong to $\mathrm{Dom}_{2} \delta$ and

$$
\begin{equation*}
\delta U=\int_{0}^{1} \dot{U}(s) d B(s) \tag{2.21}
\end{equation*}
$$

where the stochastic integral is taken in the Itô sense.
Proof. The principle of the proof is to establish (2.21) for adapted simple processes and then pass to the limit.
Step 1. For $0 \leq s<t \leq 1$, let

$$
\dot{U}(r)=\theta_{s} \mathbf{1}(s, t](r), \text { i.e. } U(r)=\theta_{s}(t \wedge r-s \wedge r),
$$

where $\theta_{s} \in \mathbb{D}_{2,1}$ and $\theta_{r}$ is $\mathcal{F}_{s}$-measurable. According to Theorem 2.8, $U$ is in $\mathrm{Dom}_{2} \delta$ and

$$
\begin{aligned}
\delta(U) & =\theta_{s} \delta(t \wedge .-s \wedge .)-\left\langle\nabla \theta_{s}, t \wedge .-s \wedge .\right\rangle_{\mathcal{H}} \\
& =\theta_{s}(B(t)-B(s))-\int_{0}^{1} \dot{\nabla}_{\tau} \theta_{s} \mathbf{1}(s, t](\tau) \tau
\end{aligned}
$$

Now recall that according to Lemma 2.5, since $\theta_{s} \in \mathcal{F}_{s}$,

$$
\dot{\nabla}_{\tau} \theta_{s}=0 \text { if } \tau>r
$$

hence

$$
\begin{equation*}
\delta(U)=\theta_{s}(B(t)-B(s))=\int_{0}^{1} \dot{U}(r) \mathrm{d} B(r) \tag{2.22}
\end{equation*}
$$

Step 2. If $\dot{U}$ is adapted, the random variable

$$
2^{n}\left(\int_{(i-1) 2^{-n}}^{i 2^{-n}} \dot{U}(r) \mathrm{d} r\right) \text { belongs to } \mathcal{F}_{i 2^{-n}}
$$

Hence, with the notations of Theorem 2.6, we have by linearity

$$
\delta\left(U_{\mathcal{D}_{n}}\right)=\int_{0}^{1} \dot{U}_{\mathcal{D}_{n}}(r) \mathrm{d} B(r)
$$

Step 3. It remains to show that we can pass to the limit in both sides of (2.21). The application $\delta$ is continuous from $\mathbb{D}_{2,1}^{a}(\mathcal{H}) \subset \mathbb{D}_{2,1}(\mathcal{H})$ into $L^{2}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$.
Hence, Theorem 2.6 entails that

$$
\delta\left(U_{\mathcal{D}_{n}}\right) \xrightarrow[n \rightarrow \infty]{L^{2}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)} \delta(U) .
$$

Furthermore, the Itô integral is an isometry hence a continuous map from $L_{a}^{2}(\mathcal{W} \times[0,1] \rightarrow \mathbf{R} ; \mu)$ into $L^{2}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$. Hence,

$$
\int_{0}^{1} \dot{U}_{\mathcal{D}_{n}}(r) \mathrm{d} B(r) \xrightarrow[n \rightarrow \infty]{L^{2}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)} \int_{0}^{1} \dot{U}(r) \mathrm{d} B(r) .
$$

The proof is thus complete.
The Itô isometry states that for $U$ adapted

$$
\mathbf{E}\left[\left(\int_{0}^{1} \dot{U}(s) \mathrm{d} B(s)\right)^{2}\right]=\mathbf{E}\left[\int_{0}^{1}|\dot{U}(s)|^{2} \mathrm{~d} s\right]
$$

One of the most elegant formula given by the Malliavin calculus is the generalization of this identity to non-adapted integrands.
Remark 2.5. If $U \in \mathbb{D}_{2,1}(\mathcal{H})$ then $\dot{\nabla} \dot{U}$ is a.s. an Hilbert-Schmidt map on $L^{2}([0,1] \times[0,1], \lambda \otimes \lambda)$. Indeed, by the definition of the norm in $\mathbb{D}_{2,1}(\mathcal{H})$,

$$
\|U\|_{\mathbb{D}_{2,1}}^{2}=\mathbf{E}\left[\|\nabla U\|_{\mathcal{H} \otimes \mathcal{H}}^{2}\right]=\mathbf{E}\left[\int_{0}^{1} \int_{0}^{1}\left|\dot{\nabla}_{r} \dot{U}(s)\right|^{2} \mathrm{~d} r \mathrm{~d} s\right] .
$$

This ensures the almost-sure finiteness of

$$
\int_{0}^{1} \int_{0}^{1}\left|\dot{\nabla}_{r} \dot{U}(s)\right|^{2} \mathrm{~d} r \mathrm{~d} s
$$

meaning that $\dot{\nabla} \dot{U}$ is Hilbert-Schmidt with probability 1.

Lemma 2.6. If $U$ belongs to $\mathbb{D}_{2,1}^{a}(\mathcal{H})$ then $\operatorname{trace}(\nabla U \circ \nabla U)=0$.
Proof. According to Lemma 1.4,

$$
\operatorname{trace}(\nabla U \circ \nabla U)=\iint_{[0,1]^{2}} \dot{\nabla}_{r} \dot{U}(s) \dot{\nabla}_{s} \dot{U}(r) \mathrm{d} r \mathrm{~d} s
$$

Since $\dot{U}(s)$ is $\mathcal{F}_{s}$-measurable, $\dot{\nabla}_{r} \dot{U}(s)=0$ if $r>s$. Similarly, $\dot{\nabla}_{s} \dot{U}(r)=0$ if $s>r$. Hence, the product is zero $\lambda \otimes \lambda$ almost-surely. It follows that the integral is null.

Theorem 2.9 ( $L^{2}$ norm of divergence). The space $\mathbb{D}_{1,2}(\mathcal{H})$ is included in $\mathrm{Dom}_{2} \delta$ and for $U \in \mathbb{D}_{1,2}(\mathcal{H})$,

$$
\begin{equation*}
\mathbf{E}\left[\delta U^{2}\right]=\mathbf{E}\left[\|U\|_{\mathcal{H}}^{2}\right]+\mathbf{E}[\operatorname{trace}(\nabla U \circ \nabla U)] \tag{2.23}
\end{equation*}
$$

Lemma 2.7. For $k \geq 1$, for $V \in \mathbb{D}_{2,1}\left(\mathcal{H}^{\otimes(k)}\right)$, for $x \in \mathcal{H}^{\otimes(k)}$, for $h \in \mathcal{H}$,

$$
\left\langle\nabla\langle V, x\rangle_{\mathcal{H} \otimes(k)}, h\right\rangle_{\mathcal{H}}=\langle\nabla V, x \otimes h\rangle_{\mathcal{H} \otimes(k+1)}
$$

Proof. For the sake of simplicity, we give the proof for $k=1$. The general case is handled similarly. Going back to the definition of the scalar product in $\mathcal{H}$, we have

$$
\left\langle\nabla\langle V, x\rangle_{\mathcal{H} \otimes(k)}, h\right\rangle_{\mathcal{H}}=\int_{0}^{1} \dot{\nabla}_{s}\left(\int_{0}^{1} \dot{V}(r) \dot{x}(r) \mathrm{d} r\right) \dot{h}(s) \mathrm{d} s
$$

Approximate the inner integral by Riemann sums and pass to the limit to show that

$$
\dot{\nabla}_{s}\left(\int_{0}^{1} \dot{V}(r) \dot{x}(r) \mathrm{d} r\right)=\int_{0}^{1} \dot{\nabla}_{s} \dot{V}(r) \dot{x}(r) \mathrm{d} r
$$

first for $V$ such that $(r, s) \longmapsto \dot{\nabla}_{s} \dot{V}(r)$ is continuous and then by density for all $V \in \mathbb{D}_{2,1}(\mathcal{H})$. Hence the result.

Lemma 2.8. For $U \in \mathbb{D}_{2,2}(\mathcal{H})$, for any $h, k, l \in \mathcal{H}$,

$$
\begin{aligned}
\left\langle\nabla^{(2)}\langle U, h\rangle_{\mathcal{H}}, k \otimes l\right\rangle_{\mathcal{H} \otimes \mathcal{H}} & =\left\langle\nabla^{(2)} U, h \otimes k \otimes l\right\rangle_{\mathcal{H} \otimes \mathcal{H}} \\
& =\left\langle\nabla^{(2)} U, h \otimes l \otimes k\right\rangle_{\mathcal{H} \otimes \mathcal{H}} .
\end{aligned}
$$

That means the map $\nabla^{(2)} U$ is a symmetric operator.
Proof. Step 1. The first equality is a consequence of Lemma 2.7.
Step 2. For $F \in \mathcal{S}, F=f\left(\delta h_{1}, \cdots, \delta h_{M}\right)$, in virtue of the Schwarz theorem for crossed derivatives of functions of several variables,

$$
\begin{aligned}
\left\langle\nabla^{(2)} F, k \otimes l\right\rangle_{\mathcal{H} \otimes \mathcal{H}} & =\sum_{i, j=1}^{n} \partial_{i j}^{2} f\left(\delta h_{1}, \cdots, \delta h_{M}\right)\left\langle k, h_{i}\right\rangle_{\mathcal{H}}\left\langle l, h_{j}\right\rangle_{\mathcal{H}} \\
& =\sum_{i, j=1}^{n} \partial_{j i}^{2} f\left(\delta h_{1}, \cdots, \delta h_{M}\right)\left\langle l, h_{j}\right\rangle_{\mathcal{H}}\left\langle k, h_{i}\right\rangle_{\mathcal{H}} \\
& =\left\langle\nabla^{(2)} F, l \otimes k\right\rangle_{\mathcal{H} \otimes \mathcal{H}}
\end{aligned}
$$

The proof is thus complete.
Proof (Proof of Theorem 2.9). For $U \in \mathbb{D}_{1,2}(\mathcal{H}), U$ takes its values in $\mathcal{H}$ so that we can write

$$
U=\sum_{n \geq 0}\left\langle U, h_{n}\right\rangle_{\mathcal{H}} h_{n},
$$

for $\left(h_{n}, n \geq 0\right)$ a complete orthonormal basis of $\mathcal{H}$. The series

$$
U_{N}=\sum_{n=0}^{N}\left\langle U, h_{n}\right\rangle_{\mathcal{H}} h_{n} \text { and } \nabla U_{N}=\sum_{n=0}^{N} \nabla\left\langle U, h_{n}\right\rangle_{\mathcal{H}} h_{n}
$$

converge in $L^{2}(\mathcal{W} \rightarrow \mathcal{H} ; \mu)$ and $\left.L^{2}(\mathcal{W} \rightarrow \mathcal{H} \otimes \mathcal{H} ; \mu)\right)$ respectively.
According to (2.17),

$$
\delta U_{N}=\sum_{n=0}^{N}\left\langle U, h_{n}\right\rangle_{\mathcal{H}} \delta h_{n}-\sum_{n=0}^{N}\left\langle\nabla U, h_{n} \otimes h_{n}\right\rangle_{\mathcal{H} \otimes \mathcal{H}} .
$$

Thus,

$$
\begin{aligned}
\nabla \delta U_{N}=\sum_{n=0}^{N}\left\{\left\langle\nabla U, h_{n}\right\rangle_{\mathcal{H}} \delta h_{n}+\left\langle U, h_{n}\right\rangle_{\mathcal{H}}\right. & h_{n} \\
& \left.-\nabla\left(\left\langle\nabla U, h_{n} \otimes h_{n}\right\rangle_{\mathcal{H} \otimes \mathcal{H}}\right)\right\}
\end{aligned}
$$

Consequently, in virtue of Lemma 2.7,

$$
\begin{aligned}
\mathbf{E}\left[\delta U_{N} \delta U_{N}\right] & =\sum_{n, k \geq 0}^{N} \mathbf{E}\left[\left\langle U, h_{k}\right\rangle_{\mathcal{H}}\left\langle\nabla U, h_{n} \otimes h_{k}\right\rangle_{\mathcal{H} \otimes \mathcal{H}} \delta h_{n}\right] \\
& +\sum_{n, k \geq 0}^{N} \mathbf{E}\left[\left\langle U, h_{n}\right\rangle_{\mathcal{H}}\left\langle U, h_{k}\right\rangle_{\mathcal{H}}\left\langle h_{n}, h_{k}\right\rangle_{\mathcal{H}}\right] \\
& -\sum_{n, k \geq 0}^{N} \mathbf{E}\left[\left\langle U, h_{k}\right\rangle_{\mathcal{H}}\left\langle\nabla^{(2)} U, h_{n} \otimes h_{n} \otimes h_{k}\right\rangle_{\mathcal{H}^{\otimes(3)}}\right] \\
& =A_{1}+A_{2}-A_{3} .
\end{aligned}
$$

On the one hand, Parseval equality yields

$$
\begin{aligned}
& A_{2}=\sum_{n, k \geq 0} \mathbf{E}\left[\left\langle U, h_{n}\right\rangle_{\mathcal{H}}\left\langle U, h_{k}\right\rangle_{\mathcal{H}}\left\langle h_{n}, h_{k}\right\rangle_{\mathcal{H}}\right] \\
&=\sum_{n \geq 0} \mathbf{E}\left[\left\langle U, h_{n}\right\rangle_{\mathcal{H}}^{2}\right]=\mathbf{E}\left[\|U\|_{\mathcal{H}}^{2}\right]
\end{aligned}
$$

Apply once more the integration by parts formula in $A_{1}$ :

$$
\begin{aligned}
A_{1} & =\sum_{n, k \geq 0}^{N} \mathbf{E}\left[\left\langle\nabla U, h_{k} \otimes h_{n}\right\rangle_{\mathcal{H} \otimes \mathcal{H}}\left\langle\nabla U, h_{n} \otimes h_{k}\right\rangle_{\mathcal{H} \otimes \mathcal{H}}\right] \\
& +\sum_{n, k \geq 0}^{N} \mathbf{E}\left[\left\langle U, h_{k}\right\rangle_{\mathcal{H}}\left\langle\nabla^{(2)} U, h_{n} \otimes h_{k} \otimes h_{n}\right\rangle_{\mathcal{H} \otimes(3)}\right] \\
& =\operatorname{trace}\left(\nabla U_{N} \circ \nabla U_{N}\right)+A_{3},
\end{aligned}
$$

since $\nabla^{(2)}$ is a symmetric operator, cf. Lemma 2.8.
Step 3. In brief, we have proved so far that

$$
\mathbf{E}\left[\delta U_{N}^{2}\right]=\left\|U_{N}\right\|_{L^{2}(W ; \mathcal{H})}^{2}+\mathbf{E}\left[\operatorname{trace}\left(\nabla U_{N} \circ \nabla U_{N}\right)\right]
$$

Then, Eqn. (1.18) entails that

$$
\begin{aligned}
\mathbf{E}\left[\delta\left(U_{N}-U_{K}\right)^{2}\right] \leq\left\|U_{N}-U_{K}\right\|_{L^{2}(\mathcal{W} \rightarrow}^{2} & \mathcal{H} ; \mu) \\
& +\left\|\nabla U_{N}-\nabla U_{K}\right\|_{L^{2}(\mathcal{W} \rightarrow \mathcal{H} \otimes \mathcal{H} ; \mu)}^{2} .
\end{aligned}
$$

Thus, the sequence $\left(\delta U_{N}, N \geq 0\right)$ is Cauchy in $L^{2}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$ thus convergent towards a limit temporarily denoted by $\eta \in L^{2}(\mathcal{W} \rightarrow \mathbf{R} ; \mu)$. Hence, for $F \in \mathbb{D}_{1,2}$,

$$
\mathbf{E}\left[\langle\nabla F, U\rangle_{\mathcal{H}}\right]=\lim _{N \rightarrow \infty} \mathbf{E}\left[\left\langle\nabla F, U_{N}\right\rangle_{\mathcal{H}}\right]=\lim _{N \rightarrow \infty} \mathbf{E}\left[F \delta U_{N}\right]=\mathbf{E}[F \eta]
$$

By the very definition of the divergence, this means that $U \in \operatorname{Dom}_{2} \delta$ and $\delta U=\eta=\lim _{N \rightarrow \infty} \delta U_{N}$.

## Exercises

Exercise 2.1. For $h \in \mathcal{H}$, show that

$$
\begin{equation*}
\nabla^{(k)} \exp \left(\delta h-\frac{1}{2}\|h\|_{\mathcal{H}}^{2}\right)=\exp \left(\delta h-\frac{1}{2}\|h\|_{\mathcal{H}}^{2}\right) h^{\otimes(k)} \tag{2.24}
\end{equation*}
$$

Exercise 2.2. Let $F \in \mathbb{D}_{p, 1}$ and $\epsilon>0$. Set $\phi_{\epsilon}(x)=\sqrt{x^{2}+\epsilon^{2}}$.

1. Show that $\phi_{\epsilon}(F) \in \mathbb{D}_{p, 1}$.
2. Show that $|F| \in \mathbb{D}_{p, 1}$ and that

$$
\dot{\nabla}_{s}|F|= \begin{cases}\dot{\nabla}_{s} F & \text { if } F>0 \\ 0 & \text { if } F=0 \\ -\dot{\nabla}_{s} F & \text { if } F<0\end{cases}
$$

3. If $G \in \mathbb{D}_{p, 1}$, compute $\nabla(F \vee G)$.

Let $B$ be the standard Brownian motion on $[0,1]$ and $M=\sup _{t \in[0,1]} B(s)$.
Let $\mathbf{Q} \cap[0,1]=\left\{t_{n}, n \geq 0\right\}$. Consider

$$
M_{n}=\sup _{s \in\left\{t_{1}, \cdots, t_{n}\right\}} B(s)
$$

We admit that $B$ attains its maximum at a unique almost-surely. Let $T=\operatorname{argmax}_{s \in[0,1]} B(s)$.
4. Show that $M_{n}$ belongs to $\mathbb{D}_{p, 1}$ and compute $\dot{\nabla} M_{n}$.
5. Prove that $M \in \mathbb{D}_{p, 1}$ and that $\nabla M=\mathbf{1}_{[0, T]}$.

Exercise 2.3 (Iterated divergence). For $U \in \mathcal{S}(\mathcal{H})$, i.e.

$$
U=\sum_{j=1}^{n} f_{j}\left(\delta h_{1}, \cdots, \delta h_{m}\right) v_{j}
$$

where $\left(v_{1}, \cdots, v_{n}\right)$ belong to $\mathcal{H}$ and $f_{j}$ in the Schwartz space on $\mathbf{R}^{m}$. Let $\delta^{(2)}$ defined by the duality

$$
\mathbf{E}\left[\delta^{2} u^{\otimes(2)} G\right]=\mathbf{E}\left[\left\langle u^{\otimes(2)}, \nabla^{(2)} G\right\rangle_{\mathcal{H} \otimes \mathcal{H}}\right]
$$

for any $G \in \mathbb{D}_{2,2}$. Show that

$$
\delta^{2}\left(U^{\otimes(2)}\right)=(\delta U)^{2}-\|U\|_{\mathcal{H}}^{2}-\operatorname{trace}(\nabla U \circ \nabla U)-2 \delta\left(\langle\nabla U, U\rangle_{\mathcal{H}}\right)
$$

Exercise 2.4 (Stratonovitch integral). The Itô integral has a major drawback: Its differential is not given by the usual formula but by the Itô formula. On the other hand, the Stratonovitch integral does satisfy the usual rule of differentiation but does not give to a martingale! We see in this exercise that the Stratonovitch integral can be computed with $\delta$ and $\nabla$. For $T_{n}=\left\{0=t_{0}<t_{1}=1 / n<\ldots<t_{n}=1\right\}$, let

$$
\mathrm{d} B_{T_{n}}(t)=\sum_{i=0}^{n-1} \frac{B\left(t_{i+1}\right)-B\left(t_{i}\right)}{t_{i+1}-t_{i}} \mathbf{1}_{\left[t_{i}, t_{i+1}\right]}(t) \mathrm{d} t
$$

be the linear affine interpolation of $B$. For any $\mathcal{H}$-valued random variable $U$, consider the Riemann-like sum

$$
S_{T_{n}}^{U}=\sum_{i=0}^{n-1} \frac{B\left(t_{i+1}\right)-B\left(t_{i}\right)}{t_{i+1}-t_{i}} \int_{t_{i}}^{t_{i+1}} \dot{U}(t) \mathrm{d} t .
$$

The process $U$ is said to be Stratonovitch integrable if the sequence $\left(S_{T_{n}}^{U}, n \geq\right.$ 0 ) converges in probability as $n$ goes to infinity.

Assume that $U$ belongs to $\mathbb{D}_{1,2}(\mathcal{H})$ and that the map

$$
\begin{aligned}
{[0,1] \times[0,1] } & \longrightarrow \mathbf{R} \\
(s, t) & \longmapsto \dot{\nabla}_{s} \dot{U}(t)
\end{aligned}
$$

is continuous.

1. Show that $U$ is Stratonovitch integrable and

$$
\lim _{n \rightarrow \infty} S_{T_{n}}^{U}=\delta U+\int_{0}^{1} \dot{\nabla}_{r} \dot{U}(r) \mathrm{d} r
$$

Indication: Verify that

$$
S_{T_{n}}^{U}=\sum_{i=0}^{n-1} \frac{1}{t_{i+1}-t_{i}} \delta\left(I^{1}\left(\mathbf{1}_{\left[t_{i}, t_{i+1}\right]}\right)\right) \int_{t_{i}}^{t_{i+1}} \dot{U}(t) d t
$$

Apply (2.17).
2. Find

$$
\lim _{n \rightarrow \infty} \sum_{i=0}^{n-1} \frac{1}{2}\left(\dot{U}\left(t_{i}\right)+\dot{U}\left(t_{i+1}\right)\right)\left(B\left(t_{i+1}\right)-B\left(t_{i}\right)\right)
$$

