# INTEGRATION OF BROWNIAN VECTOR FIELDS 

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#### Abstract

Using the Wiener chaos decomposition, we show that strong solutions of non-Lipschitzian stochastic differential equations are given by random Markovian kernels. The example of Sobolev flows is studied in some detail, exhibiting interesting phase transitions.


0. Introduction. The purpose of this paper is to present an extended notion of strong solutions of stochastic differential equations (SDEs) driven by Wiener processes. These solutions can be defined on rather general spaces, in the context of Dirichlet forms.

More interesting, they are not always given by flows of maps but by flows of Markovian kernels, which means splitting can occur. Coalescent flows also appear as solutions of these SDEs. Conditions are given under which coalescence and splitting occur or not.

A variety of examples are studied. The case of isotropic Sobolev flows on the sphere or on the Euclidean space shows in particular that splitting is related to hyperinstability and coalescence to hyperstability. These notions (which will be developed in Sections 9 and 10) are related to the explosion of the Lyapunov exponent toward $+\infty$ and $-\infty$.

The typical example we have in mind is the Brownian motion on a Riemannian manifold. We consider a covariance on vector fields which induces the Riemannian metric on each tangent space. When the covariance function has enough regularity, it is known that one can solve the linear SDE driven by the canonical Wiener process associated with this covariance [or in other terms to the local characteristics associated with this covariance (see Section 3 below)] and get a multiplicative Brownian motion on the diffeomorphism group, which moves every point as a Brownian motion (see [18] or [23]). However, models related to turbulence theory produce natural examples where the regularity condition is not satisfied. Except for the work of Darling [7], where strong solutions are not considered, these SDEs have not been really studied. The idea is to define the solutions by their Wiener chaos expansion in terms of the heat semigroup. We call it the statistical solution. A similar expansion was given by Krylov and Veretennikov in [17], for SDEs with strong solutions.

In this form, they appear as a semigroup of operators, and the fact that these operators are Markovian is not clearly visible in the formula. To prove this, we

[^0]consider an independent realization of the Brownian motion on the manifold and couple it with the given Wiener process on vector fields using certain martingales. Then the Markovian random operators which constitute the strong solution are obtained by filtering the Brownian motion with respect to this Wiener process. They determine the law of a canonical weak solution of the equation given the Wiener process on vector fields. This construction has been adequately generalized to be presented in the case of symmetric diffusions on a locally compact metric space. It is a convenient and well-studied framework but this assumption could clearly be relaxed (in particular to the framework of coercive forms). Relations with particle representations and filtering of stochastic partial differential equations (SPDEs) can be observed (see [19]).

The example of Sobolev flows is studied in detail on Euclidean spaces and spheres and is of major interest especially in dimensions 2 and 3 where an interesting phase diagram is given in terms of the two parameters determining the Sobolev norm on vector fields: The differentiability index and the relative weight of gradients and divergence free fields (compressibility).

Some of these results have been given in the note [21], and a preliminary version of this work was released in [22]. They are directly connected and partially motivated by a series of works of Gawedzki and Kupiainen on turbulent advection $[2,13,14]$.

1. Covariance function on a manifold. Let $X$ be a manifold. A covariance function $C$ on $T^{*} X$ is a symmetric map from $T^{*} X^{2}$ in $\mathbb{R}$ such that, for any $(x, y) \in X^{2}, C$ restricted to $T_{x}^{*} X \times T_{y}^{*} X$ is bilinear and such that, for any $n$-tuples $\left(\xi_{1}, \ldots, \xi_{n}\right)$ of $T^{*} X$,

$$
\begin{equation*}
\sum_{i, j} C\left(\xi_{i}, \xi_{j}\right) \geq 0 \tag{1.1}
\end{equation*}
$$

For any $\xi=(x, u) \in T^{*} X$, let $C_{\xi}$ be the vector field such that, for any $\xi^{\prime}=$ $(y, v) \in T^{*} X$,

$$
\left\langle C_{\xi}(y), v\right\rangle=C\left(\xi, \xi^{\prime}\right)
$$

Let $H_{0}$ be the vector space generated by the vector fields $C_{\xi}$. Let us define the bilinear form on $H_{0},\langle\cdot, \cdot\rangle_{H}$ such that

$$
\begin{equation*}
\left\langle C_{\xi}, C_{\xi^{\prime}}\right\rangle_{H}=C\left(\xi, \xi^{\prime}\right) \tag{1.2}
\end{equation*}
$$

As (1.1) is satisfied, the bilinear form $\langle\cdot, \cdot\rangle_{H}$ is a scalar product on $H_{0}$. We denote by $\|\cdot\|_{H}$ the norm associated with $\langle\cdot, \cdot\rangle_{H}$.

Let $H$ be the separate completion of $H_{0}$ with respect to $\|\cdot\|_{H} \cdot\left(H,\langle\cdot, \cdot\rangle_{H}\right)$ is a separable Hilbert space and we will designate it as the self-reproducing space associated with the covariance function $C . H$ is constituted of vector fields on $X$ and, for any $h \in H$ and any $\xi=(x, u) \in T^{*} X$,

$$
\begin{equation*}
\left\langle C_{\xi}, h\right\rangle_{H}=\langle h(x), u\rangle \tag{1.3}
\end{equation*}
$$

Let $\left(e_{k}\right)_{k}$ be an orthonormal basis of $H$, then (1.3) implies that, for any $\xi=(x, u) \in T^{*} X$,

$$
\begin{equation*}
C_{\xi}=\sum_{k}\left\langle e_{k}(x), u\right\rangle e_{k} \tag{1.4}
\end{equation*}
$$

This equation implies that, for any $\xi^{\prime}=(y, v) \in T^{*} X$,

$$
\begin{equation*}
C\left(\xi, \xi^{\prime}\right)=\sum_{k}\left\langle e_{k}(x), u\right\rangle\left\langle e_{k}(y), v\right\rangle \tag{1.5}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
C=\sum_{k} e_{k} \otimes e_{k} \tag{1.6}
\end{equation*}
$$

REMARK 1.1. On the other hand, if we start with a countable family of vector fields $\left(V_{k}\right)_{k}$, such that for any $\xi=(x, u) \in T^{*} X, \sum_{k}\left\langle V_{k}(x), u\right\rangle^{2}<\infty$, it is possible to define a covariance function on $X$ by the formula

$$
C=\sum_{k} V_{k} \otimes V_{k}
$$

Examples of isotropic covariances are given in Sections 9 and 10. See also [1].
Now assume a Riemannian metric $\langle\cdot, \cdot\rangle_{x}$ is given on $X$, the linear bundles $T X$ and $T^{*} X$ can be identified. We will now suppose that the covariance is bounded by the metric, that is, that

$$
C(\xi, \xi) \leq\langle u, u\rangle_{x}
$$

for any $\xi=(x, u) \in T^{*} X$. Note that this condition implies that $|h(x)|_{x} \leq\|h\|_{H}$ for any $h \in H$.

Let us denote by $m(d x)$ the volume element on $X$. Given any differentiable function $f$ such that $|\nabla f|$ is square integrable, we can map it linearly into $D f$ in the Hilbert tensor product $L^{2}(m) \hat{\otimes} H$ setting $\langle D f, g \otimes h\rangle=\int_{X} g(x) \times$ $\langle\nabla f(x), h(x)\rangle_{x} m(d x)$ for all $g \in L^{2}(m)$ and $h \in H$.

Note that

$$
\begin{equation*}
\|D f\|_{H}^{2}(x) \leq|\nabla f(x)|^{2} \tag{1.7}
\end{equation*}
$$

[this notation comes from the identification $L^{2}(m) \hat{\otimes} H$ with the $L^{2}$-space of $H$-valued functions on $X$ ] and that

$$
\begin{equation*}
\|D f\|_{L^{2}(m) \otimes H}^{2} \leq \int|\nabla f|^{2} d m \tag{1.8}
\end{equation*}
$$

2. Covariance function bounded by a Dirichlet form. We can extend these notions to the framework of local Dirichlet forms. Let $X$ be a locally compact separable metric space and let $m$ be a positive Radon measure on $X$ such that $\operatorname{Supp}[m]=X$.

Let $(\mathcal{E}, \mathcal{F})$ be a regular Dirichlet space, $\mathcal{F} \subset L^{2}(X, m)$. We suppose that the Dirichlet form is local and conservative. We denote by $P_{t}$ the associated Markovian semigroup, by $A$ its generator and by $\mathscr{D}(A)$ the domain of $A$. We also suppose that $m$ is an invariant measure, hence that $P_{t} 1=1$. We also assume that, for any $f \in \mathcal{F}_{b}=L^{\infty}(m) \cap \mathcal{F}$, there exists $\Gamma(f, f) \in L^{1}(m)$ such that, for any $g \in \mathcal{F}_{b}$,

$$
\begin{equation*}
2 \mathcal{E}(f g, f)-\mathcal{E}\left(f^{2}, g\right)=\int g \Gamma(f, f) d m \tag{2.1}
\end{equation*}
$$

$\Gamma$ can be extended to $\mathcal{F}$ and we denote by $\Gamma(f, g)$ the $L^{1}(m)$-valued bilinear form on $\mathcal{F}^{2}$, where for any $(f, g) \in \mathcal{F}^{2}, \Gamma(f, g)=\frac{1}{4}(\Gamma(f+g, f+g)-\Gamma(f-g$, $f-g)$ ). A sufficient condition for the existence of $\Gamma$ (see [3], Corollary 4.2.3) is that $\mathscr{D}(A)$ contains a subspace $E$ of $\mathscr{D}(A)$, dense in $\mathcal{F}$, such that

$$
\forall f \in E, \quad f^{2} \in \mathscr{D}(A)
$$

Then, for $(f, g) \in E^{2}$,

$$
\begin{equation*}
\Gamma(f, g)=A(f g)-f A g-g A f \tag{2.2}
\end{equation*}
$$

A necessary and sufficient condition for the existence of the energy density (or carré du champ operator) $\Gamma$ is given in Theorem 4.2.2 in [3].

FUNDAMENTAL EXAMPLE 2.1. $X$ is a Riemannian manifold with metric $\langle\cdot, \cdot\rangle, m$ is the volume measure, $\mathcal{F}=H^{1}(X)$ and, for any $(f, g) \in \mathcal{F}^{2}$,

$$
\mathcal{E}(f, g)=\frac{1}{2} \int_{X}\langle\nabla f, \nabla g\rangle d m
$$

In this case, $\Gamma(f, g)=\langle\nabla f, \nabla g\rangle$.
Let $H$ be a separable Hilbert space and $D$ a linear map from $\mathcal{F}$ into the Hilbert tensor product $L^{2}(m) \hat{\otimes} H$ such that, for any $f \in \mathcal{F}$,

$$
\begin{equation*}
\|D f(x)\|_{H}^{2} \leq \Gamma(f, f)(x) \tag{2.3}
\end{equation*}
$$

$m(d x)$-a.e. The most interesting case is when there is equality in (2.3).
We define the covariance function $C$ as a bilinear map from $\mathcal{F} \times \mathcal{F}$ into $L^{2}(m \otimes m)$ by

$$
\begin{equation*}
\langle C(f, g), u \otimes v\rangle_{L^{2}(m \otimes m)}=\int_{X^{2}}\langle D f(x), D g(y)\rangle_{H} u(x) v(y) m(d x) m(d y) \tag{2.4}
\end{equation*}
$$

Note that

$$
\begin{equation*}
\langle C(f, f), u \otimes u\rangle_{L^{2}(m \otimes m)} \leq 2 \mathcal{E}(f, f)\|u\|_{L^{2}(m)}^{2} \tag{2.5}
\end{equation*}
$$

We say that $C$ is a covariance function bounded by the Dirichlet form $(\mathcal{E}, \mathcal{F})$.

REMARK 2.2. Alternatively, we could define the covariance as a positive symmetric bilinear map from $\mathcal{F} \times \mathcal{F}$ in $L^{2}(m \otimes m)$ [i.e., such that, for any $u_{i} \in L^{2}(m)$ and any $f_{i} \in \mathcal{F}$,

$$
\begin{equation*}
\left.\int \sum_{i, j} u_{i} \otimes u_{j} C\left(f_{i}, f_{j}\right) d m^{\otimes 2} \geq 0\right] \tag{2.6}
\end{equation*}
$$

such that

$$
\begin{equation*}
\langle C(f, f), u \otimes u\rangle_{L^{2}(m \otimes m)} \leq 2 \mathscr{E}(f, f)\|u\|_{L^{2}(m)}^{2} \tag{2.7}
\end{equation*}
$$

and construct as before a Hilbert space such that (2.3) holds.
Indeed, we define $H$ as the separated closure of the space $H_{0}$ spanned by elements of the form $u \otimes f$, with $u \otimes f \in L^{2}(m) \otimes \mathcal{F}$, and equipped with the (possibly degenerate) scalar product

$$
\langle u \otimes f, v \otimes g\rangle_{H}=\langle C(f, g), u \otimes v\rangle_{L^{2}(m \otimes m)}
$$

And for $f \in \mathcal{F}, D f$ is defined such that for any $u \otimes v \otimes g \in L^{2}(m) \otimes L^{2}(m) \otimes \mathcal{F}$,

$$
\langle D f, u \otimes v \otimes g\rangle_{L^{2}(m) \otimes} \hat{\otimes}=\langle C(f, g), u \otimes v\rangle_{L^{2}(m \otimes m)}
$$

For any $h \in H$ and $f \in \mathcal{F}$ define $D_{h} f=\langle D f, h\rangle_{H}$, which belongs to $L^{2}(m)$. Then, for any orthonormal basis $\left(e_{k}\right)_{k}$ of $H$,

$$
\begin{equation*}
C=\sum_{k} D_{e_{k}} \otimes D_{e_{k}} . \tag{2.8}
\end{equation*}
$$

Moreover, for any $f \in \mathcal{F}$,

$$
\begin{equation*}
\|D f\|_{H}^{2}=\sum_{k}\left(D_{e_{k}} f\right)^{2} . \tag{2.9}
\end{equation*}
$$

Note that condition (2.3) implies that, for any finite family $\left(u_{i}, f_{i}\right) \in$ $L^{\infty}(m) \times \mathcal{F}$,

$$
\begin{equation*}
\sum_{i, j} u_{i} u_{j} D\left(f_{i}, f_{j}\right) \leq \sum_{i, j} u_{i} u_{j} \Gamma\left(f_{i}, f_{j}\right) \tag{2.10}
\end{equation*}
$$

where $D(f, g)$ denotes $\langle D f(x), D g(x)\rangle_{H}=\sum_{k} D_{e_{k}} f(x) D_{e_{k}} g(x)$. When the $u_{i}$ are step functions with discontinuities in a set of zero measure, (2.10) is satisfied as $\sum_{i, j} u_{i} u_{j} D\left(f_{i}, f_{j}\right)=\left|D\left(\sum_{i} u_{i} f_{i}\right)\right|^{2}$. Then we can extend this result to any family $\left(u_{i}\right)$ by density in $L^{2}(\Gamma(f, f) d m)$ for every $f \in \mathcal{F}$.

REMARK 2.3. It is clear that, given a covariance $C$ on $T^{*} X$ as in Section 1, we can build the self-reproducing space $H$ consisting of vector fields and the mapping $D: H^{1}(X) \rightarrow L^{2}(m) \hat{\otimes} H$ so as to construct a covariance function as in Section 2. Now suppose conversely that we have a separable Hilbert space $H$, a linear map $D$ and a covariance $C$ as in Section 2, and suppose we are in the Riemannian case.

The condition $\|D f(x)\|_{H}^{2} \leq \Gamma(f, f)(x)=|\nabla f(x)|^{2}$ implies that $C(f, g)(x, y)$ depends only on $\nabla f(x)$ and $\nabla g(y)$, and so there is a covariance $\tilde{C}$ say on $T^{*} X$ so that $C(f, g)(x, y)=\tilde{C}(\nabla f(x), \nabla g(y))$. So in the Riemannian case, any Section 2 covariance function reduces to a Section 1 covariance function.

Further, we can now assume without any loss of generality that the separable Hilbert space $H$ is the self-reproducing space corresponding to $\tilde{C}$ and thus consists of vector fields.

REMARK 2.4. The bilinear mapping $D$ is a derivation: for any $h \in H$ and any $f \in \mathcal{F}$ such that $f^{2} \in \mathcal{F}$,

$$
\begin{equation*}
D_{h} f^{2}=2 f D_{h} f \tag{2.11}
\end{equation*}
$$

Note that in the Riemanian manifold case (Fundamental example 2.1), $D_{h} f=$ $\nabla_{h} f$ when $\Gamma=D$.

Proof. We first make the remark that

$$
\sum_{k}\left(D_{e_{k}} f^{2}-2 f D_{e_{k}} f\right)^{2}=D\left(f^{2}, f^{2}\right)-4 f D\left(f^{2}, f\right)+4 f^{2} D(f, f)
$$

Integrating this relation with respect to $m$ and using (2.10), we get that

$$
\begin{aligned}
\int \sum_{k}\left(D_{e_{k}} f^{2}-2 f D_{e_{k}} f\right)^{2} d m & \leq \int\left(\Gamma\left(f^{2}, f^{2}\right)-4 f \Gamma\left(f^{2}, f\right)+4 f^{2} \Gamma(f, f)\right) d m \\
& =0
\end{aligned}
$$

This implies that, for every $k, D_{e_{k}} f^{2}-2 f D_{e_{k}} f=0$.
3. Construction of the statistical solutions. In Fundamental example 2.1, when $X$ is a Riemannian manifold, $C$ is smooth and, when equality holds in (2.3), it is well known (see $[18,23]$ ) that a stochastic flow of diffeomorphisms on $X$ can be associated with $C$. Then, with the notation of Definition 2.1 in [23], the local characteristics of the flow are $(A, L)$, where $A=C$ and $L$ is the Laplacian on $X$.

The object of this section is to show that, in the general situation considered above, it is always possible to define a flow of Markovian kernels associated with $C$ and $(\mathcal{E}, \mathcal{F})$ (which is induced by the stochastic flow when $C$ is smooth).

Let a covariance function $C$ bounded by a Dirichlet form $(\mathcal{E}, \mathcal{F})$ on a locally compact separable metric space be given as in the preceding section [(2.3) is satisfied]. Let $W_{t}$ be a cylindrical Brownian motion on $H$ defined on some probability space $(\Omega, \mathcal{A}, P)$, that is, a Gaussian process indexed by $H \times \mathbb{R}^{+}$ with covariance matrix $\operatorname{cov}\left(W_{t}(h), W_{s}\left(h^{\prime}\right)\right)=s \wedge t\left\langle h, h^{\prime}\right\rangle_{H}$. Set $W_{t}^{k}=W_{t}\left(e_{k}\right)$. ( $W_{t}^{k} ; k \in \mathbb{N}$ ) is a sequence of independent Wiener processes and we can represent $W_{t}$ by $\sum_{k} W_{t}^{k} e_{k}$. Informally, the law of $W_{t}$ is given by

$$
\frac{1}{Z} e^{-(1 / 2) \int_{0}^{\infty}\left\|\dot{W}_{t}\right\|_{H}^{2} d t} D W
$$

Let $\mathcal{F}_{t}=\sigma\left(W_{s}^{k} ; k \in \mathbb{N} ; s \leq t\right)=\sigma\left(W_{s} ; s \leq t\right)$.

Proposition 3.1. Let $S_{t}^{0}=P_{t}$. We can define a sequence $S_{t}^{n}$ of random operators on $L^{2}(m)$ such that $E\left[\left(S_{t}^{n} f\right)^{2}\right] \leq P_{t} f^{2}$ in $L^{1}(m)$ and $S_{t}^{n}$ is $\mathcal{F}_{t}$-measurable, by the recurrence formula, in $L^{2}(m \otimes P)$ [i.e., in the Hilbert tensor product $\left.L^{2}(m) \hat{\otimes} L^{2}(P)\right]$

$$
\begin{equation*}
S_{t}^{n+1} f=P_{t} f+\sum_{k} \int_{0}^{t}\left(S_{s}^{n}\left(D_{e_{k}} P_{t-s} f\right)\right) d W_{s}^{k} \tag{3.1}
\end{equation*}
$$

REMARK. The stochastic integral in (3.1) here makes sense as a Hilbertvalued Itô integral. Recall that given a real Wiener process $W_{t}$ and a Hilbert space $H$, for any $F$ progressively measurable in $L^{2}\left(P_{W} \otimes d t\right) \hat{\otimes} H$ and any $h \in H,\left\langle\int_{s}^{t} F(u) d W_{u}, h\right\rangle_{H}=\int_{s}^{t}\langle F(u), h\rangle_{H} d W_{u}$ and $E\left[\left\|\int_{s}^{t} F(u) d W_{u}\right\|_{H}^{2}\right]=$ $\int_{s}^{t}\|F(u)\|_{H}^{2} d u$.

Proof of Proposition 3.1. Suppose we are given $S_{t}^{n}$, an $\mathcal{F}_{t}$-measurable random operator on $L^{2}(m)$ such that $E\left[\left(S_{t}^{n} f\right)^{2}\right] \leq P_{t} f^{2}$.

Let $f \in L^{2}(m)$. For any positive $t, P_{t} f \in \mathcal{F}$ and $D_{e_{k}} P_{t-s} f$ is well defined and belongs to $L^{2}(m)$,

$$
\begin{aligned}
E\left[\left(S_{t}^{n+1} f\right)^{2}\right] & =\left(P_{t} f\right)^{2}+\sum_{k} \int_{0}^{t} E\left[\left(S_{s}^{n}\left(D_{e_{k}} P_{t-s} f\right)\right)^{2}\right] d s, \quad m \text {-a.e. } \\
& \leq\left(P_{t} f\right)^{2}+\int_{0}^{t} P_{s}\left(\left|D P_{t-s} f\right|^{2}\right) d s \\
& \leq\left(P_{t} f\right)^{2}+\int_{0}^{t} P_{s}\left(\Gamma\left(P_{t-s} f, P_{t-s} f\right)\right) d s
\end{aligned}
$$

For $f \in L^{\infty}(m) \cap L^{2}(m), \frac{\partial}{\partial s} P_{s}\left(\left(P_{t-s} f\right)^{2}\right)=P_{s}\left(\Gamma\left(P_{t-s} f, P_{t-s} f\right)\right)$ and

$$
\begin{equation*}
P_{t} f^{2}=\left(P_{t} f\right)^{2}+\int_{0}^{t} P_{s}\left(\Gamma\left(P_{t-s} f, P_{t-s} f\right)\right) d s \tag{3.2}
\end{equation*}
$$

An approximation by truncation shows that (3.2) remains true for $f \in L^{2}(m)$ and $E\left[\left(S_{t}^{n+1} f\right)^{2}\right] \leq P_{t} f^{2}$.

REMARK. The definition of $S_{t}^{n}$ is independent of the choice of the basis on $H$.

In the following, we use the canonical realization of the processes $W_{t}^{k}$. They are defined as coordinate functions on $\Omega=C\left(\mathbb{R}^{+}, \mathbb{R}\right)^{\mathbb{N}}$, with the product Wiener measure $P$. We denote by $\theta_{t}$ the natural shift on $\Omega$ such that $W_{t+s}^{k}-W_{t}^{k}=W_{s}^{k} \circ \theta_{t}$.

Recall that an operator on $L^{2}(m)$ is called Markovian if and only if it preserves positivity and maps 1 into 1 (or, more precisely, if $m$ is not finite, its natural extension to positive functions maps 1 into 1 ).

THEOREM 3.2. The family of random operators $S_{t}^{n}$ converges in $L^{2}(P)$ toward a one-parameter family of $\mathscr{F}_{t}$-adapted Markovian operators $S_{t}$ such that the following hold:
(a) $S_{t+s}=S_{t}\left(S_{s} \circ \theta_{t}\right)$, for any $s, t \geq 0$;
(b) $\forall f \in L^{2}(m), S_{t} f$ is uniformly continuous with respect to $t$ in $L^{2}(m \otimes P)$;
(c) $E\left[\left(S_{t} f\right)^{2}\right] \leq P_{t} f^{2}$, for any $f \in L^{2}(m)$;
(d) $S_{t} f=P_{t} f+\sum_{k} \int_{0}^{t} S_{s}\left(D_{e_{k}} P_{t-s} f\right) d W_{s}^{k}$, for any $f \in L^{2}(m)$;
(e) $S_{t} f=f+\sum_{k} \int_{0}^{t} S_{s}\left(D_{e_{k}} f\right) d W_{s}^{k}+\int_{0}^{t} S_{s}(A f) d s$, for any $f \in \mathscr{D}(A)$.
$S_{t}$ is uniquely characterized by (c) and (d) or by (a), (c) and (e). When $\Gamma=D$, we call it the statistical solution of the SDE [see (3.22) below]

$$
\begin{equation*}
\forall f \in \mathscr{D}(A), \quad d f\left(X_{t}\right)=\sum_{k} D_{e_{k}} f\left(X_{t}\right) d W_{t}^{k}+A f\left(X_{t}\right) d t . \tag{3.3}
\end{equation*}
$$

Note that this SDE does not always have a strong solution in the usual sense.

Proof. The convergence of $S_{t}^{n}$ is immediate since, for any $n \geq 1, J_{t}^{n} f=$ $S_{t}^{n} f-S_{t}^{n-1} f$ is in the Hilbert tensor product of the $n$th Wiener chaos of $L^{2}(P)$ with $L^{2}(m), S_{t} f=P_{t} f+\sum_{n=1}^{\infty} J_{t}^{n} f$ and $\left(P_{t} f\right)^{2}+\sum_{n \geq 1} E\left[\left(J_{t}^{n} f\right)^{2}\right]=$ $\lim _{n \rightarrow \infty} E\left[\left(S_{t}^{n} f\right)^{2}\right] \leq P_{t} f^{2}$. It is clear that $S_{t}$ is $\mathcal{F}_{t}$-adapted and satisfies (c). Part (d) is obtained by taking the limit in the recurrence formula of the proposition.

Since

$$
J_{t}^{n} f=\sum_{k_{1}, \ldots, k_{n}} \int_{0<s_{1}<\cdots<s_{n}<t} P_{s_{1}} D_{e_{k_{1}}} P_{s_{2}-s_{1}} \cdots D_{e_{k_{n}}} P_{t-s_{n}} f d W_{s_{1}}^{k_{1}} \cdots d W_{s_{n}}^{k_{n}}
$$

we have $J_{t+s}^{n}=\sum_{k \leq n} J_{t}^{k}\left(J_{s}^{n-k} \circ \theta_{t}\right)$ (the $k$ th term corresponds to

$$
\begin{gathered}
\sum_{k_{1}, \ldots, k_{n}} \int_{0<s_{1}<\cdots<s_{k}<s<s_{k+1}<\cdots<s_{n}<t+s} P_{s_{1}} D_{e_{k_{1}}} P_{s_{2}-s_{1}} \cdots D_{e_{k_{n}}} P_{t+s-s_{n}} f \\
\left.\times d W_{s_{1}}^{k_{1}} \cdots d W_{s_{n}}^{k_{n}}\right)
\end{gathered}
$$

We deduce (a) from this relation.
The uniqueness of a solution of (d) satisfying (c) follows directly from the uniqueness of the Wiener chaos decomposition, obtained by iteration of (d): Let $T_{t}$ designate another solution of (d) and (c). Then, for any $f \in L^{2}(m)$ and any integer $n$,
$T_{t} f=S_{t}^{n-1} f+\sum_{k_{1}, \ldots, k_{n}} \int_{0<s_{1}<\cdots<s_{n}<t} T_{s_{1}} D_{e_{k_{1}}} P_{s_{2}-s_{1}} \cdots D_{e_{k_{n}}} P_{t-s_{n}} f d W_{s_{1}}^{k_{1}} \cdots d W_{s_{n}}^{k_{n}}$.

The second term of the right-hand side of the preceding equation is orthogonal to the first one since its integrands are $L^{2}$. Indeed,

$$
\begin{aligned}
& \sum_{k_{1}, \ldots, k_{n}} E\left[\int_{0<s_{1}<\cdots<s_{n}<t}\left(T_{s_{1}} D_{e_{k_{1}}} P_{s_{2}-s_{1}} \cdots D_{e_{k_{n}}} P_{t-s_{n}} f\right)^{2} d s_{1} \cdots d s_{n}\right] \\
& \quad \leq \sum_{k_{1}, \ldots, k_{n}} \int_{0<s_{1}<\cdots<s_{n}<t} P_{s_{1}}\left(\left|D_{e_{k_{1}}} P_{s_{2}-s_{1}} \cdots D_{e_{k_{n}}} P_{t-s_{n}} f\right|^{2}\right) d s_{1} \cdots d s_{n} \\
& \quad \leq \sum_{k_{2}, \ldots, k_{n}} \int_{0<s_{2}<\cdots<s_{n}<t} P_{s_{2}}\left(\left|D_{e_{k_{2}}} P_{s_{3}-s_{2}} \cdots D_{e_{k_{n}}} P_{t-s_{n}} f\right|^{2}\right) d s_{2} \cdots d s_{n}
\end{aligned}
$$

using (2.3) and (3.2) and by induction is smaller than $P_{t} f^{2}$.
This proves that the Wiener chaos decompositions of $T_{t} f$ and $S_{t} f$ are the same and therefore $T_{t}=S_{t}$.

Proof of (b). Let us remark that, for any positive $\varepsilon, S_{t+\varepsilon}-S_{t}=S_{t}\left(S_{\varepsilon} \circ\right.$ $\left.\theta_{t}-I\right)$. As $S_{t}$ and $S_{\varepsilon} \circ \theta_{t}$ are independent and $m$ is invariant under $P_{t}$, for any $f \in L^{2}(m)$,

$$
\begin{align*}
\int E\left[\left(S_{t+\varepsilon} f-S_{t} f\right)^{2}\right] d m & \leq \int E\left[P_{t}\left(S_{\varepsilon} \circ \theta_{t} f-f\right)^{2}\right] d m \\
& \leq \int E\left[\left(S_{\varepsilon} \circ \theta_{t} f-f\right)^{2}\right] d m  \tag{3.4}\\
& \leq \int\left(P_{\varepsilon} f^{2}-2 f P_{\varepsilon} f+f^{2}\right) d m \\
& \leq 2\|f\|_{L^{2}(m)}\left\|f-P_{\varepsilon} f\right\|_{L^{2}(m)} .
\end{align*}
$$

Therefore, $\lim _{\varepsilon \rightarrow 0}\left\|S_{t+\varepsilon} f-S_{t} f\right\|_{L^{2}(m \otimes P)}=0$, uniformly in $t$.
Remark 3.3. Note also the convergence in $L^{2}(m \otimes P)$ of $P_{\varepsilon} S_{t} f$ toward $S_{t} f$ when $\varepsilon \rightarrow 0$. Indeed $\left\|P_{\varepsilon} S_{t} f-S_{t} f\right\|_{L^{2}(m \otimes P)}^{2}=E\left[\left\|P_{\varepsilon} S_{t} f-S_{t} f\right\|_{L^{2}(m)}^{2}\right]$ and $\left\|P_{\varepsilon} S_{t} f-S_{t} f\right\|_{L^{2}(m)}^{2}$ converges toward 0 when $\varepsilon$ goes to 0 and is dominated by $4\left\|S_{t} f\right\|_{L^{2}(m)}^{2}$.

Proof of (e). Let us remark that, for any $\varepsilon$ and $t$ positive,

$$
\begin{align*}
S_{t+\varepsilon} f-S_{t} f & =S_{t}\left(P_{\varepsilon} f+\sum_{k} \int_{0}^{\varepsilon} S_{u} \circ \theta_{t}\left(D_{e_{k}} P_{\varepsilon-u} f\right) d W_{u}^{k} \circ \theta_{t}-f\right)  \tag{3.5}\\
& =S_{t}\left(P_{\varepsilon} f-f\right)+\sum_{k} \int_{t}^{t+\varepsilon} S_{s}\left(D_{e_{k}} P_{t+\varepsilon-s} f\right) d W_{s}^{k}
\end{align*}
$$

Hence using (3.5) for $t=\frac{i}{n} t$ and $\varepsilon=\frac{1}{n} t$, for $f \in \mathscr{D}(A)$,

$$
\begin{aligned}
& S_{t} f-f-\sum_{k} \int_{0}^{t} S_{s}\left(D_{e_{k}} f\right) d W_{s}^{k}-\int_{0}^{t} S_{s}(A f) d s \\
&=\sum_{i=0}^{n-1}\left[S_{(i / n) t}\left(P_{t / n} f-f\right)+\sum_{k} \int_{(i / n) t}^{((i+1) / n) t} S_{s}\left(D_{e_{k}} P_{((i+1) / n) t-s} f\right) d W_{s}^{k}\right. \\
&\left.\quad-\int_{(i / n) t}^{((i+1) / n) t} S_{s}(A f) d s-\sum_{k} \int_{(i / n) t}^{((i+1) / n) t} S_{s}\left(D_{e_{k}} f\right) d W_{s}^{k}\right]
\end{aligned}
$$

$$
=A_{1}(n)+A_{2}(n)+A_{3}(n),
$$

with

$$
\begin{align*}
& A_{1}(n)=\sum_{i=0}^{n-1} S_{(i / n) t}\left(P_{t / n} f-f-\frac{t}{n} A f\right)  \tag{3.6}\\
& A_{2}(n)=\sum_{i=0}^{n-1} \int_{(i / n) t}^{((i+1) / n) t}\left(S_{(i / n) t}(A f)-S_{s}(A f)\right) d s  \tag{3.7}\\
& A_{3}(n)=\sum_{i=0}^{n-1} \sum_{k} \int_{(i / n) t}^{((i+1) / n) t} S_{s}\left(D_{e_{k}}\left(P_{((i+1) / n) t-s} f-f\right)\right) d W_{s}^{k} . \tag{3.8}
\end{align*}
$$

First, using the fact that $m$ is $P_{t}$-invariant,

$$
\begin{equation*}
\left\|A_{1}(n)\right\|_{L^{2}(m \otimes P)} \leq n\left\|P_{t / n} f-f-\frac{t}{n} A f\right\|_{L^{2}(m)}=o(1) . \tag{3.9}
\end{equation*}
$$

After, we remark that

$$
\begin{aligned}
& \left\|\int_{(i / n) t}^{((i+1) / n) t}\left(S_{(i / n) t}(A f)-S_{s}(A f)\right) d s\right\|_{L^{2}(m \otimes P)}^{2} \\
& \quad \leq \frac{t}{n} \int_{(i / n) t}^{((i+1) / n) t}\left\|S_{(i / n) t}(A f)-S_{s}(A f)\right\|_{L^{2}(m \otimes P)}^{2} d s .
\end{aligned}
$$

As $S_{t}(A f)$ is uniformly continuous in $L^{2}(m \otimes P)$, there exists $\varepsilon(x)$ such that $\lim _{x \rightarrow 0} \varepsilon(x)=0$ and $\left\|S_{(i / n) t}(A f)-S_{s}(A f)\right\|_{L^{2}(m \otimes P)}^{2} \leq \varepsilon\left(\frac{t}{n}\right)$ for any $s \in\left[\frac{i}{n} t\right.$, $\left.\frac{i+1}{n} t\right]$. Hence we get

$$
\left\|\int_{(i / n) t}^{((i+1) / n) t}\left(S_{(i / n) t}(A f)-S_{s}(A f)\right) d s\right\|_{L^{2}(m \otimes P)}^{2} \leq \frac{t^{2}}{n^{2}} \varepsilon\left(\frac{t}{n}\right)
$$

and $\left\|A_{2}(n)\right\|_{L^{2}(m \otimes P)}=o(1)$.

At last, as the different terms in the sum in (3.8) are orthogonal,
$\left\|A_{3}(n)\right\|_{L^{2}(m \otimes P)}^{2}$

$$
\begin{aligned}
& =\sum_{i=0}^{n-1} \sum_{k} \int E\left[\left(\int_{(i / n) t}^{((i+1) / n) t} S_{s}\left(D_{e_{k}}\left(P_{((i+1) / n) t-s} f-f\right)\right) d W_{s}^{k}\right)^{2}\right] d m \\
& \leq \sum_{i=0}^{n-1} \iint_{(i / n) t}^{((i+1) / n) t}\left|D\left(P_{((i+1) / n) t-s} f-f\right)\right|^{2} d s d m \\
& \leq n \int_{0}^{t / n} \int\left|D\left(P_{s} f-f\right)\right|^{2} d m d s \\
& \leq n \int_{0}^{t / n} \mathcal{E}\left(P_{s} f-f, P_{s} f-f\right) d s
\end{aligned}
$$

As $\lim _{s \rightarrow 0} \mathcal{E}\left(P_{s} f-f, P_{s} f-f\right)=0,\left\|A_{3}(n)\right\|_{L^{2}(m \otimes P)}=o(1)$.
Taking the limit as $n$ goes to $\infty$, this shows that $\| S_{t} f-f-\sum_{k} \int_{0}^{t} S_{s}\left(D_{e_{k}} f\right) \times$ $d W_{s}^{k}-\int_{0}^{t} S_{s}(A f) d s \|_{L^{2}(m \otimes P)}=0$.

Proof that (a), (c) AND (e) IMPLY (d). Take $f \in L^{2}(m)$ and $\varepsilon$ positive, assuming (e),

$$
\begin{aligned}
& S_{t} P_{\varepsilon} f-P_{t} P_{\varepsilon} f-\sum_{k} \int_{0}^{t} S_{s}\left(D_{e_{k}} P_{t-s} P_{\varepsilon} f\right) d W_{s}^{k} \\
& =\sum_{i=0}^{n-1}\left[S_{((i+1) / n) t}\left(P_{t-((i+1) / n) t} P_{\varepsilon} f\right)-S_{(i / n) t}\left(P_{t-(i / n) t} P_{\varepsilon} f\right)\right. \\
& \left.\quad-\sum_{k} \int_{(i / n) t}^{((i+1) / n) t} S_{s}\left(D_{e_{k}}\left(P_{t-s} P_{\varepsilon} f\right)\right) d W_{s}^{k}\right] \\
& = \\
& =B_{1}(n)+B_{2}(n)+B_{3}(n),
\end{aligned}
$$

with

$$
\begin{gather*}
B_{1}(n)=\sum_{i=0}^{n-1} \sum_{k} \int_{(i / n) t}^{((i+1) / n) t} S_{s}\left(D_{e_{k}}\left(P_{t-((i+1) / n) t} P_{\varepsilon} f-P_{t-s} P_{\varepsilon} f\right)\right) d W_{s}^{k}  \tag{3.10}\\
B_{2}(n)=-\sum_{i=0}^{n-1} S_{(i / n) t}\left(P_{t-(i / n) t} P_{\varepsilon} f-P_{t-((i+1) / n) t} P_{\varepsilon} f\right.  \tag{3.11}\\
\left.-\frac{t}{n} A P_{t-((i+1) / n) t} P_{\varepsilon} f\right) ;
\end{gather*}
$$

$$
\begin{equation*}
B_{3}(n)=\sum_{i=0}^{n-1} \int_{(i / n) t}^{((i+1) / n) t}\left(S_{s}-S_{(i / n) t}\right)\left(A P_{t-((i+1) / n) t} P_{\varepsilon} f\right) d s, \tag{3.12}
\end{equation*}
$$

since

$$
\begin{aligned}
& S_{((i+1) / n) t}\left(P_{t-((i+1) / n) t} P_{\varepsilon} f\right) \\
& =S_{(i / n) t}\left(P_{t-((i+1) / n) t} P_{\varepsilon} f\right)+\int_{(i / n) t}^{((i+1) / n) t} S_{S}\left(A P_{t-((i+1) / n) t} P_{\varepsilon} f\right) d s \\
& \quad+\sum_{k} \int_{(i / n) t}^{((i+1) / n) t} S_{S} D_{e_{k}} P_{t-((i+1) / n) t} P_{\varepsilon} f d W_{s}^{k}
\end{aligned}
$$

Since the different terms in the sum in (3.10) are orthogonal,

$$
\begin{aligned}
& \left\|B_{1}(n)\right\|_{L^{2}(m \otimes P) \leq}^{2} \leq \sum_{i=0}^{n-1} \iint_{(i / n) t}^{((i+1) / n) t}\left|D\left(P_{t-((i+1) / n) t} P_{\varepsilon} f-P_{t-s} P_{\varepsilon} f\right)\right|^{2} d s d m \\
& \leq
\end{aligned}
$$

as $\mathcal{E}\left(P_{t} f, P_{t} f\right) \leq \mathcal{E}(f, f)$ for any positive $t$ and any $f \in L^{2}(m)$.
Equation (3.13) implies that $\left\|B_{1}(n)\right\|_{L^{2}(m \otimes P)}=o(1)$ [as $\lim _{s \rightarrow 0} \mathscr{E}\left(P_{S} P_{\varepsilon} f-\right.$ $\left.\left.P_{\varepsilon} f, P_{S} P_{\varepsilon} f-P_{\varepsilon} f\right)=0\right]$;

$$
\begin{aligned}
& \left\|B_{2}(n)\right\|_{L^{2}(m \otimes P)} \\
& \qquad \begin{array}{l}
\leq \sum_{i=0}^{n-1} \| S_{(i / n) t}\left(P_{t-(i / n) t} P_{\varepsilon} f-P_{t-((i+1) / n) t} P_{\varepsilon} f\right. \\
\left.\quad-\frac{t}{n} A P_{t-((i+1) / n) t} P_{\varepsilon} f\right) \|_{L^{2}(m \otimes P)} \\
\leq \\
\leq \sum_{i=0}^{n-1}\left(\int\left(P_{((i+1) / n) t} P_{\varepsilon} f-P_{(i / n) t} P_{\varepsilon} f-\frac{t}{n} A P_{(i / n) t} P_{\varepsilon} f\right)^{2} d m\right)^{1 / 2} \\
\leq n\left\|P_{t / n} P_{\varepsilon} f-P_{\varepsilon} f-\frac{t}{n} A P_{\varepsilon} f\right\|_{L^{2}(m)},
\end{array}
\end{aligned}
$$

hence, $\left\|B_{2}(n)\right\|_{L^{2}(m \otimes P)}=o(1)$.
Note that if $Q_{t} f=E\left[S_{t} f\right]$, (e) implies that, for any $f \in \mathscr{D}(A)$,

$$
Q_{t} f=f+\int_{0}^{t} Q_{s}(A f) d s
$$

Then $\frac{\partial}{\partial_{s}} Q_{s} P_{t-s} f=0$ for any $f \in L^{2}(m)$ and $0<s<t\left[\right.$ then $\left.P_{t-s} f \in \mathscr{D}(A)\right]$ and we have $Q_{t} f=P_{t} f$. With this remark and the fact that (a) and (c) are satisfied, we see that (b) and (3.4) are satisfied [see the proof of (b)]. Using (3.4), we have

$$
\begin{aligned}
\|\left(S_{s}-\right. & \left.S_{(i / n) t}\right)\left(A P_{t-((i+1) / n) t} P_{\varepsilon} f\right) \|_{L^{2}(m \otimes P)}^{2} \\
\leq & 2\left\|A P_{t-((i+1) / n) t} P_{\varepsilon} f\right\|_{L^{2}(m)} \\
& \times\left\|A P_{t-((i+1) / n) t} P_{\varepsilon} f-P_{s-(i / n) t} A P_{t-((i+1) / n) t} P_{\varepsilon} f\right\|_{L^{2}(m)} \\
\leq & 2\left\|A P_{\varepsilon} f\right\|_{L^{2}(m)}\left\|A P_{\varepsilon} f-P_{s-(i / n) t} A P_{\varepsilon} f\right\|_{L^{2}(m)} \\
\leq & 4\left\|A P_{\varepsilon} f\right\|_{L^{2}(m)}^{2} .
\end{aligned}
$$

Hence,

$$
\begin{aligned}
& \left\|B_{3}(n)\right\|_{L^{2}(m \otimes P)}^{2} \\
& \quad \leq \sum_{i=0}^{n-1} \frac{t}{n} \int_{(i / n) t}^{((i+1) / n) t}\left\|\left(S_{s}-S_{(i / n) t}\right)\left(A P_{t-((i+1) / n) t} P_{\varepsilon} f\right)\right\|_{L^{2}(m \otimes P)}^{2} d s \\
& \quad \leq \frac{4 t^{2}}{n}\left\|A P_{\varepsilon} f\right\|_{L^{2}(m)}^{2}
\end{aligned}
$$

Taking the limit as $n$ goes to $\infty$, this shows that (d) is satisfied for $P_{\varepsilon} f$, with $f \in L^{2}(m)$ and $\varepsilon$ positive.

At last, since $\left\|S_{t} P_{\varepsilon} f-S_{t} f\right\|_{L^{2}(m \otimes P)} \leq\left\|P_{\varepsilon} f-f\right\|_{L^{2}(m)}$ [because (c) is satisfied], $\left\|P_{t+\varepsilon} f-P_{t} f\right\|_{L^{2}(m \otimes P)} \leq\left\|P_{\varepsilon} f-f\right\|_{L^{2}(m)}$ and $\| \sum_{k} \int_{0}^{t} S_{s}\left(D_{e_{k}} P_{t-s}\left(P_{\varepsilon} f-\right.\right.$ $f)) d W_{s}^{k} \|_{L^{2}(m \otimes P)}^{2} \leq t \mathcal{E}\left(P_{\varepsilon} f-f, P_{\varepsilon} f-f\right)$. Taking the limit as $\varepsilon$ goes to 0 , we prove that (d) is satisfied for any $f \in L^{2}(m)$.

Proof that $S_{t}$ is Markovian. A more concise proof of this fact has been given in [21], relying on Wiener exponentials and Girsanov formula. The advantage of the following proof is to be more explanatory, to give a relation with weak solutions and to yield a construction of the process law associated with the statistical solution $S_{t}$.

Let $\left(\Omega^{\prime}, \mathcal{G}, \mathcal{q}_{t}, X_{t}, P_{x}\right)$ be a Hunt process associated with $(\mathcal{E}, \mathcal{F})$ (see [11]); we take a canonical version with $\Omega^{\prime}=C\left(\mathbb{R}^{+}, X\right)$. Let $W_{t}=\sum_{k} W_{t}^{k} e_{k}$ be a cylindrical Brownian motion on $H$, independent of the Markov process $X_{t}$.

Let $\mathcal{M}$ be the space of the martingales additive functionals, $g_{t}$-adapted such that if $M \in \mathcal{M}, E_{x}\left[M_{t}^{2}\right]<\infty, E_{x}\left[M_{t}\right]=0$ q.e. and $e(M)<\infty$, where $e(M)=$ $\sup _{t>0} \frac{1}{2 t} E_{m}\left[M_{t}^{2}\right]\left[\right.$ with $\left.P_{m}=\int P_{x} d m(x)\right]$. ( $\left.\mathcal{M}, e\right)$ is a Hilbert space (see [11]).

For $f \in \mathcal{F}, M^{f} \in \mathcal{M}$ denotes the martingale part of the semimartingale $f\left(X_{t}\right)-$ $f\left(X_{0}\right)$. For $g \in C_{K}(X) \subset L^{2}(\Gamma(f, f) d m)\left[C_{K}(X)\right.$ designates the space of functions continuous with compact support], we denote by $g . M^{f} \in \mathcal{M}$ the martingale $\int_{0}^{t} g\left(X_{s}\right) d M_{s}^{f}$. Then $\mathcal{M}_{0}=\left\{\sum_{i=1}^{n} g_{i} . M^{f_{i}} ; n \in \mathbb{N}, g_{i} \in C_{K}(X), f_{i} \in \mathcal{F}\right\}$ is dense
in $\mathcal{M}$ (see [11], Lemma 5.6.3), and $e\left(\sum_{i} g_{i} . M^{f_{i}}\right)=\frac{1}{2} \sum_{i, j} \int g_{i} g_{j} \Gamma\left(f_{i}, f_{j}\right) d m$ (see [11], Theorem 5.2.3 and 5.6.1).

LEMmA 3.4. For every $(M, N) \in \mathcal{M} \times \mathcal{M}$, there exists $\Gamma(M, N) \in L^{1}(m)$ such that

$$
\begin{equation*}
\langle M, N\rangle_{t}=\int_{0}^{t} \Gamma(M, N)\left(X_{s}\right) d s \tag{3.14}
\end{equation*}
$$

where $\langle\cdot, \cdot\rangle_{t}$ is the usual martingale bracket. For $(f, g) \in \mathcal{F}, \Gamma\left(M^{f}, M^{g}\right)=$ $\Gamma(f, g)$.

Note that Lemma 3.4 implies that $e(M, N)=\frac{1}{2} \int \Gamma(M, N) d m$.
In Fundamental example 2.1, $X_{t}$ is the Brownian motion on $X, M_{t}^{f}$ is the Itô integral $\int_{0}^{t}\left\langle d f\left(X_{s}\right), d X_{s}\right\rangle, \Gamma$ is the inverse Riemannian metric and $\mathcal{M}$ can be identified with the space of 1 -forms equipped with the $L^{2}$-norm associated with the metric.

Proof of Lemma 3.4. When $f \in \mathscr{F}$, it follows from Theorem 5.2.3 in [11] that

$$
\left\langle M^{f}, M^{f}\right\rangle_{t}=\int_{0}^{t} \Gamma(f, f)\left(X_{s}\right) d s
$$

For $M=\sum_{i} h_{i} \cdot M^{f_{i}}, N=\sum_{j} k_{j} \cdot M^{g_{j}}$, two martingales of $\mathcal{M}_{0}$,

$$
\begin{equation*}
\langle M, N\rangle_{t}=\sum_{i, j} \int_{0}^{t} h_{i} k_{j} \Gamma\left(f_{i}, g_{j}\right)\left(X_{s}\right) d s=\int_{0}^{t} \Gamma(M, N)\left(X_{s}\right) d s \tag{3.15}
\end{equation*}
$$

with $\Gamma(M, N)=\sum_{i, j} h_{i} k_{j} \Gamma\left(f_{i}, g_{j}\right) . \Gamma$ is a bilinear mapping from $\mathcal{M}_{0} \times \mathcal{M}_{0}$ in $L^{1}(m) . \Gamma$ is continuous since, for any $(M, N) \in \mathcal{M}_{0} \times \mathcal{M}_{0}$,

$$
\begin{aligned}
\int|\Gamma(M, N)| d m & \leq \int \Gamma(M, M)^{1 / 2} \Gamma(N, N)^{1 / 2} d m \\
& \leq 2 e(M)^{1 / 2} e(N)^{1 / 2}
\end{aligned}
$$

It follows that $\Gamma$ can be extended to $\mathcal{M} \times \mathcal{M}$.
Take $M \in \mathcal{M}$ and an approximating sequence $M_{n} \in \mathcal{M}_{0}$. Then $e\left(M_{n}-M\right)$ converges toward $0, M_{n}$ converges toward $M$ in $L^{2}\left(P_{x}\right)$ and $\left\langle M_{n}, M_{n}\right\rangle_{t}$ converge in $L^{1}\left(P_{x}\right)$ toward $\langle M, M\rangle_{t}$ for almost every $x$ (see [11], Section 5-2). This proves that $\langle M, M\rangle_{t}=\int_{0}^{t} \Gamma(M, M)\left(X_{s}\right) d s$.

LEMMA 3.5. Suppose $m$ is bounded. Then for any $h \in H$, there exists a unique continuous martingale in $\mathcal{M}, N^{h}$ such that, for any $f \in \mathcal{F}, e\left(N^{h}, M^{f}\right)=$ $\frac{1}{2} \int D_{h} f d m$ and $\frac{d}{d t}\left\langle N^{h}, M^{f}\right\rangle_{t}=D_{h} f\left(X_{t}\right)$. In addition, $e\left(N^{h}\right) \leq \frac{1}{2} m(X)\|h\|^{2}$ and $\left\langle N^{h}\right\rangle_{t} \leq\|h\|^{2} t$.

In the Riemannian manifold case (Fundamental example 2.1), $N_{t}^{h}=\int_{0}^{t}\left\langle h\left(X_{s}\right)\right.$, $\left.d X_{s}\right\rangle$ when $\Gamma=D$.

Proof of Lemma 3.5. For $h=\sum_{k} \lambda_{k} e_{k} \in H$, let us define a linear form $\alpha_{h}$ on $\mathcal{M}_{0}$ such that, for any $M=\sum_{i=1}^{n} g_{i} . M^{f_{i}} \in \mathcal{M}_{0}, \alpha_{h}(M)=\frac{1}{2} \sum_{i=1}^{n} \int g_{i} D_{h} f_{i} d m$;

$$
\begin{aligned}
\left(\alpha_{h}(M)\right)^{2} & =\left(\sum_{k} \lambda_{k} \frac{1}{2} \sum_{i=1}^{n} \int g_{i} D_{e_{k}} f_{i} d m\right)^{2} \\
& \leq \frac{1}{4}\|h\|^{2} m(X) \sum_{i, j} \int g_{i} g_{j} D\left(f_{i}, f_{j}\right) d m \leq \frac{1}{2}\|h\|^{2} m(X) e(M)
\end{aligned}
$$

This proves that $\alpha_{h}$ is continuous on $\mathcal{M}_{0}$ and can be extended to a continuous linear form on $\mathcal{M}$ such that $\alpha_{h}(M) \leq \frac{1}{\sqrt{2}}\|h\| \sqrt{m(X) e(M)}$. With this form is associated a unique $N^{h} \in \mathcal{M}$ such that $\alpha_{h}(M)=e\left(N^{h}, M\right)$.

Note that, for any $g \in C_{K}(X)$ and $f \in \mathcal{F}$, we have $\int g D_{h} f d m=\int \Gamma\left(N^{h}\right.$, $\left.g . M^{f}\right) d m=\int g \Gamma\left(N^{h}, M^{f}\right) d m$. This is satisfied for every $g \in C_{K}(X)$; therefore, for any $f \in \mathcal{F}, \Gamma\left(N^{h}, M^{f}\right)=D_{h} f$.

Note that we also have, for $M \in \mathcal{M}_{0}$,

$$
\Gamma\left(N^{h}, M\right) \leq\|h\| \Gamma(M, M)^{1 / 2}
$$

which implies that $\left\langle N^{h}\right\rangle_{t} \leq\|h\|^{2} t$.
REMARK 3.6. When $m$ is not bounded, $N^{h}$ can be defined as a local martingale such that, for any compact $K$ and any $f \in \mathcal{F}, \mathbb{1}_{K} \cdot N^{h} \in \mathcal{M}$, $e\left(\mathbb{1}_{K} \cdot N^{h}, M^{f}\right)=\frac{1}{2} \int_{K} D_{h} f d m$. In addition, $e\left(\mathbb{1}_{K} \cdot N^{h}\right) \leq \frac{1}{2} m(K)\|h\|^{2}$.

Let $\gamma_{k l}$ be a function on $X$ such that $\frac{d}{d t}\left\langle N^{e_{l}}, N^{e_{k}}\right\rangle_{t}=\gamma_{k l}\left(X_{t}\right)$. Lemma 3.5 implies that the matrix $A=\left(\left(\delta_{k l}-\gamma_{k l}\right)\right)$ is positive (as $\frac{d}{d t}\left\langle N^{h}\right\rangle_{t} \leq\|h\|^{2}$ ). Therefore, it is possible to find a matrix $R$ such that $R^{2}=A$.

REMARK 3.7. If for any $f \in \mathcal{F}, \Gamma(f, f)=\|D f\|_{H}^{2}$, then for any $f \in \mathcal{F}_{b}$,

$$
\begin{equation*}
M_{t}^{f}=\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d N_{s}^{e_{k}} \tag{3.16}
\end{equation*}
$$

$D_{e_{k}} f=\sum_{l} D_{e_{l}} f \gamma_{k l}\left(X_{t}\right)$ and the positive symmetric matrix $P=\left(\left(\gamma_{k l}\right)\right)$ is a projector. In this case, $R=I-P$.

Proof. Set $Q_{t}^{f}=\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d N_{s}^{e_{k}}, Q^{f} \in \mathcal{M}$. Then, for any $M=$ $\sum_{i=1}^{n} g_{i} \cdot M^{f_{i}} \in \mathcal{M}_{0}$,

$$
\begin{aligned}
\left\langle Q^{f}, M\right\rangle_{t} & =\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d\left\langle N^{e_{k}}, M\right\rangle_{s} \\
& =\sum_{k} \sum_{i=1}^{n} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) g_{i}\left(X_{s}\right) D_{e_{k}} f_{i}\left(X_{s}\right) d s \\
& =\sum_{i=1}^{n} \int_{0}^{t} g_{i}\left(X_{s}\right) D\left(f, f_{i}\right)\left(X_{s}\right) d s=\left\langle M^{f}, M\right\rangle_{t} .
\end{aligned}
$$

This implies that, for any $M \in \mathcal{M}, e\left(Q^{f}, M\right)=e\left(M^{f}, M\right)$ and $Q^{f}=M^{f}$.
Since, by Lemma 3.5, $\frac{d}{d t}\left\langle M^{f}, N^{e_{k}}\right\rangle_{t}=D_{e_{k}} f\left(X_{t}\right)$, we get that

$$
D_{e_{k}} f\left(X_{t}\right)=\frac{d}{d t}\left\langle Q^{f}, N^{e_{k}}\right\rangle_{t}=\sum_{l} D_{e_{l}} f\left(X_{t}\right) \frac{d}{d t}\left\langle N^{e_{l}}, N^{e_{k}}\right\rangle_{t}=\left(\sum_{l} D_{e_{l}} f \gamma_{k l}\right)\left(X_{t}\right) .
$$

This relation implies that $N_{t}^{e_{k}}=\sum_{l} \int_{0}^{t} \gamma_{k l}\left(X_{s}\right) d N_{s}^{e_{l}}$ (this is easy to check, considering $\frac{d}{d t}\left\langle N^{e_{k}}, M\right\rangle_{t}$ with $\left.M \in \mathcal{M}_{0}\right)$. From this, we see that $\gamma_{k l}=\sum_{i} \gamma_{k i} \gamma_{i l}$ (i.e., $P^{2}=P$ ).

Set $\widetilde{W}_{t}^{k}=N_{t}^{e_{k}}+\sum_{l} \int_{0}^{t} R_{k l}\left(X_{s}\right) d W_{s}^{l}$ and $\widetilde{W}_{t}=\sum_{k} \widetilde{W}_{t}^{k} e_{k}$.
In the Riemannian manifold case, when $\Gamma(f, f)=\|D f\|_{H}^{2}$ for any $f \in \mathcal{F}$, denoting $C_{\xi}$ by $C_{(x, u)}$ when $u \in T_{x} X$ and $\xi=(x, u)$ we have

$$
d \widetilde{W}_{t}=d W_{t}+C_{\left(X_{t}, d X_{t}\right)}-C_{\left(X_{t}, d W_{t}\left(X_{t}\right)\right)}
$$

and

$$
d \widetilde{W}_{t}^{k}=d W_{t}^{k}+\left\langle e_{k}\left(X_{t}\right), d X_{t}\right\rangle-\sum_{l}\left\langle e_{k}\left(X_{t}\right), e_{l}\left(X_{t}\right)\right\rangle d W_{t}^{l} .
$$

In this case, $R$ is a projector (see the remark above).

## Lemma 3.8. $\quad\left(\widetilde{W}_{t}^{k}\right)_{k}$ is a sequence of independent Brownian motion.

Proof. Since $\widetilde{W}_{t}^{k}$ is a continuous martingale, we just have to compute $\frac{d}{d t}\left\langle\widetilde{W}_{t}^{k}, \widetilde{W}_{t}^{l}\right\rangle_{t}$ :

$$
\frac{d}{d t}\left\langle\widetilde{W}_{t}^{k}, \widetilde{W}_{t}^{l}\right\rangle_{t}=\gamma_{k l}+R_{k l}^{2}=\delta_{k l} .
$$

This implies the lemma.
Let $\mu$ be an initial distribution of the form $h m$, with $h$ a positive function in $L^{2}(m) \cap L^{1}(m)$ and for $f \in L^{2}(m)$ define $\widetilde{S}_{t} f$ by the conditional expectation

$$
\begin{equation*}
\widetilde{S}_{t} f\left(X_{0}\right)=E_{\mu}\left[f\left(X_{t}\right) \mid \sigma\left(X_{0}, \widetilde{W}_{s}^{k} ; k \in \mathbb{N} ; s \leq t\right)\right] . \tag{3.17}
\end{equation*}
$$

(One can check easily that this definition does not depend on $h$.) Remark that, as $X_{t}$ is Markovian and $W_{t}$ has independent increments,

$$
\begin{equation*}
\widetilde{S}_{t} f\left(X_{0}\right)=E_{\mu}\left[f\left(X_{t}\right) \mid \sigma\left(X_{0}, \widetilde{W}_{s}^{k} ; k \in \mathbb{N} ; s \geq 0\right)\right] \tag{3.18}
\end{equation*}
$$

In the same way, we see that $\widetilde{S}_{t}$ satisfies the multiplicative cocycle property (a).
LEMMA 3.9. For any $f \in \mathscr{D}(A)$ and $\mu$ an initial distribution absolutely continuous with respect to $m$,

$$
\widetilde{S}_{t} f=f+\sum_{k} \int_{0}^{t} \widetilde{S}_{s}\left(D_{e_{k}} f\right) d \widetilde{W}_{s}^{k}+\int_{0}^{t} \widetilde{S}_{s}(A f) d s, \quad P_{\mu}-a . s
$$

Proof. For any $f \in \mathscr{D}(A)$, we have

$$
\begin{equation*}
f\left(X_{t}\right)=f\left(X_{0}\right)+M_{t}^{f}+\int_{0}^{t} A f\left(X_{s}\right) d s \tag{3.19}
\end{equation*}
$$

It is clear that $E .\left[\int_{0}^{t} A f\left(X_{s}\right) d s \mid \sigma\left(\widetilde{W}_{s}^{k} ; \quad k \in \mathbb{N} ; \quad s \leq t\right)\right]=\int_{0}^{t} \widetilde{S}_{s} A f\left(X_{s}\right) d s$, as (3.17) is satisfied. Let $Z_{t}=\sum_{k} \int_{0}^{t} H_{s}^{k} d \widetilde{W}_{s}^{k} \in L^{2}\left(\sigma\left(\widetilde{W}_{s}^{k} ; k \in \mathbb{N} ; s \leq t\right)\right)$,

$$
\begin{aligned}
E .\left[Z_{t} M_{t}^{f}\right] & =\sum_{k} E \cdot\left[\int_{0}^{t} H_{s}^{k} d\left\langle\widetilde{W}^{k}, M^{f}\right\rangle_{s}\right] \\
& =\sum_{k} E \cdot\left[\int_{0}^{t} H_{s}^{k} D_{e_{k}} f\left(X_{s}\right) d s\right] \\
& =\sum_{k} E \cdot\left[\int_{0}^{t} H_{s}^{k} \widetilde{S}_{s}\left(D_{e_{k}} f\right) d s\right] \\
& =E \cdot\left[Z_{t} \sum_{k} \int_{0}^{t} \widetilde{S}_{s}\left(D_{e_{k}} f\right) d \widetilde{W}_{s}^{k}\right]
\end{aligned}
$$

This proves that $E .\left[M_{t}^{f} \mid \sigma\left(\widetilde{W}_{s}^{k} ; k \in \mathbb{N} ; s \leq t\right)\right]=\sum_{k} \int_{0}^{t} \widetilde{S}_{s} D_{e_{k}} f d \widetilde{W}_{s}^{k}$.
Now, using uniqueness in Theorem 3.2 and the isomorphism $j$ between $L^{2}\left(\sigma\left(\widetilde{W}_{t}^{k} ; t \geq 0 ; k \in \mathbb{N}\right)\right)$ and $L^{2}\left(\sigma\left(W_{t}^{k} ; t \geq 0 ; k \in \mathbb{N}\right)\right)$, we see that $j \widetilde{S}_{t}=S_{t}$, which implies that $S_{t}$ is Markovian.

PROPOSITION 3.10. For any $f \in \mathcal{F}_{b}$, the martingale

$$
\begin{equation*}
P_{t}^{f}=M_{t}^{f}-\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d \widetilde{W}_{s}^{k} \tag{3.20}
\end{equation*}
$$

is orthogonal to the family of martingales $\left\{\widetilde{W}_{t}^{k} ; k \in \mathbb{N}\right\}$, in the sense of the martingale bracket (i.e., for any $k,\left\langle P^{f}, \widetilde{W}^{k}\right\rangle$. $=0$ ). For any $(f, g) \in \mathcal{F}_{b}^{2}$,

$$
\begin{equation*}
\left\langle P^{f}, P^{g}\right\rangle_{t}=\int_{0}^{t}\left(\Gamma(f, g)\left(X_{s}\right)-D(f, g)\left(X_{s}\right)\right) d s \tag{3.21}
\end{equation*}
$$

Proof. We just have to show that $\left\langle P^{f}, \widetilde{W}^{k}\right\rangle_{t}=0$ for every $f \in \mathcal{F}_{b}$ and every $k \in \mathbb{N}$, which is true as

$$
\left\langle M^{f}, \widetilde{W}^{k}\right\rangle_{t}=\left\langle M^{f}, N^{e_{k}}\right\rangle_{t}=\int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d s
$$

Let $(f, g) \in \mathcal{F}_{b}^{2}$. Then

$$
\begin{aligned}
\left\langle P^{f}, P^{g}\right\rangle_{t} & =\left\langle P^{f}, M^{g}\right\rangle_{t} \\
& =\left\langle M^{f}, M^{g}\right\rangle_{t}-\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d\left\langle\widetilde{W}^{k}, M^{g}\right\rangle_{s} \\
& =\int_{0}^{t} \Gamma(f, g)\left(X_{s}\right) d s-\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{S}\right) D_{e_{k}} g\left(X_{s}\right) d s \\
& =\int_{0}^{t} \Gamma(f, g)\left(X_{s}\right) d s-\int_{0}^{t} D(f, g)\left(X_{S}\right) d s
\end{aligned}
$$

REMARK 3.11. In the case $\Gamma(f, f)=\|D f\|_{H}^{2}$ for any $f \in \mathcal{F}$, Proposition 3.10 implies that $P_{t}^{f}=0$ and that

$$
M_{t}^{f}=\sum_{k} \int_{0}^{t} D_{e_{k}} f\left(X_{s}\right) d \widetilde{W}_{s}^{k}
$$

From this, we see that the diffusion $X_{t}$ satisfies the SDE

$$
\begin{equation*}
f\left(X_{t}\right)-f\left(X_{0}\right)=\sum_{l} \int_{0}^{t} D_{e_{l}} f\left(X_{s}\right) d \widetilde{W}_{s}^{k}+\int_{0}^{t} A f\left(X_{s}\right) d s \tag{3.22}
\end{equation*}
$$

for every $f \in \mathscr{D}(A)$. Therefore $\left(X_{t}, \widetilde{W}_{t}\right)$ appears as a weak solution of this SDE and $\widetilde{S}_{t}$ is defined by filtering $X_{t}$ with respect to $\widetilde{W}_{t}$.

Let $P_{x, \widetilde{\omega}}\left(d \omega^{\prime}\right)$ be the conditional law of the diffusion $X_{t}$, given $X_{0}$ and $\left\{\widetilde{W}_{t}\right.$; $\left.t \in \mathbb{R}^{+}\right\}$(it is independent of the choice of the initial distribution). Using the identity in law between $W$ and $\widetilde{W}$, we get a family of conditional probabilities $P_{x, \omega}\left(d \omega^{\prime}\right)$ on $C\left(\mathbb{R}^{+}, X\right)$ defined $m \otimes P$-a.e.

Remark that [with $X_{t}\left(\omega^{\prime}\right)=\omega^{\prime}(t)$ ]

$$
\begin{equation*}
S_{t} f(x, \omega)=\int f\left(X_{t}\left(\omega^{\prime}\right)\right) P_{x, \omega}\left(d \omega^{\prime}\right), \quad m \otimes P \text {-a.s. } \tag{3.23}
\end{equation*}
$$

Under $P_{x, \omega}\left(d \omega^{\prime}\right) P(d \omega), X_{t}\left(\omega^{\prime}\right)$ satisfies the $\operatorname{SDE}$ (3.3). It is a canonical weak solution of the $\operatorname{SDE}$ (3.3) on a canonical extension of the probability space on which $W$ is defined. $S_{t}$ is obtained by filtering $X_{t}$ with respect to $W$.
4. The $n$-point motion. Let $P_{t}^{(n)}$ be the family of operators on $L^{\infty}\left(m^{\otimes n}\right)$ such that, for any $\left(f_{i}\right)_{1 \leq i \leq n} \in L^{\infty}(m)$,

$$
\begin{equation*}
P_{t}^{(n)} f_{1} \otimes \cdots \otimes f_{n}=E\left[S_{t} f_{1} \otimes \cdots \otimes S_{t} f_{n}\right] \tag{4.1}
\end{equation*}
$$

$P_{t}^{(n)}$ is a Markovian semigroup on $L^{\infty}\left(m^{\otimes n}\right)$ as $S_{t}$ is Markovian and satisfies Theorem 3.2. It is easy to check that $P_{t}^{(2)}$ maps tensor products of $L^{2}(m)$ functions into $L^{2}\left(m^{\otimes 2}\right)$.

Proposition 4.1. For any family of probability laws on $X$ absolutely continuous with respect to $m,\left(\mu_{i} ; 1 \leq i \leq n\right)$,

$$
\begin{equation*}
P_{\mu_{1} \otimes \cdots \otimes \mu_{n}}^{(n)}\left(d \omega_{1}^{\prime}, \ldots, d \omega_{n}^{\prime}\right)=\int_{\Omega} P(d \omega) \bigotimes_{i=1}^{n} P_{\mu_{i}, \omega}\left(d \omega_{i}^{\prime}\right) \tag{4.2}
\end{equation*}
$$

defines a Markov process on $X^{n}$ (with initial distribution $\otimes_{i=1}^{n} \mu_{i}$ ) associated with $P_{t}^{(n)}$. We call this Markov process on $X^{n}$ the $n$-point motion.

Proof. For every family of functions in $L^{\infty}(m),\left(f_{i}\right)_{1 \leq i \leq n}, m^{\otimes n} \otimes P$-a.e. [with $\left.X_{t}^{i}\left(\omega_{i}^{\prime}\right)=\omega_{i}^{\prime}(t)\right]$,

$$
\begin{align*}
S_{t}^{\otimes n} f_{1} \otimes \cdots \otimes f_{n}\left(x_{1}, \ldots, x_{n}, \omega\right) & =\prod_{i=1}^{n} S_{t} f_{i}\left(x_{i}, \omega\right)  \tag{4.3}\\
& =\int \prod_{i=1}^{n} f_{i}\left(X_{t}^{i}\left(\omega_{i}^{\prime}\right)\right) \bigotimes_{i=1}^{n} P_{x_{i}, \omega}\left(d \omega_{i}^{\prime}\right)
\end{align*}
$$

We get the result by integrating both members of (4.3) with respect to $P(d \omega)$.

Let $D^{(n)}$ be the linear map from $H \times \mathcal{F}^{\otimes n}$ in $L^{2}\left(m^{\otimes n}\right)$ such that, for any $\left(f_{i}\right)_{1 \leq i \leq n} \in \mathcal{F}$ and $h \in H$,

$$
\begin{equation*}
D_{h}^{(n)} f_{1} \otimes \cdots \otimes f_{n}=\sum_{i=1}^{n} f_{1} \otimes \cdots \otimes D_{h} f_{i} \otimes \cdots \otimes f_{n} \tag{4.4}
\end{equation*}
$$

Proposition 4.2. For any $\left(f_{i}\right)_{1 \leq i \leq n} \in \mathscr{D}(A) \cap L^{\infty}(m)$,

$$
\begin{aligned}
S_{t}^{\otimes n} f_{1} \otimes \cdots \otimes f_{n}= & f_{1} \otimes \cdots \otimes f_{n}+\sum_{k} \int_{0}^{t} S_{s}^{\otimes n}\left(D_{e_{k}}^{(n)} f_{1} \otimes \cdots \otimes f_{n}\right) d W_{s}^{k} \\
& +\int_{0}^{t} S_{s}^{\otimes n}\left(A^{(n)} f_{1} \otimes \cdots \otimes f_{n}\right) d s
\end{aligned}
$$

where

$$
\begin{aligned}
A^{(n)} f_{1} \otimes \cdots \otimes f_{n}= & \sum_{i=1}^{n} f_{1} \otimes \cdots \otimes A f_{i} \otimes \cdots \otimes f_{n} \\
& +\sum_{1 \leq i<j \leq n} \sum_{k} f_{1} \otimes \cdots \otimes D_{e_{k}} f_{i} \otimes \cdots \otimes D_{e_{k}} f_{j} \otimes \cdots \otimes f_{n}
\end{aligned}
$$

REMARK 4.3. (a) For $n=2$, the formula extends to functions in $\mathscr{D}(A)$ and $A^{(2)} f \otimes g=A f \otimes g+f \otimes A g+C(f, g)$, where $(f, g) \in(\mathscr{D}(A))^{2}$.
(b) Taking the expectation, we see that $A^{(n)}$ is the infintesimal generator of $P_{t}^{(n)}$ on $\left(\mathscr{D}(A) \cap L^{\infty}(m)\right)^{\otimes n}$.
(c) The formula extends to $C_{K}^{2}\left(X^{n}\right)$ in the Riemannian manifold case (using for example the uniform density of sums of product functions and the regularizing effect of $P_{\varepsilon}^{\otimes n}$ ).

Proof of Proposition 4.2. This is just a straightforward application of Itô's formula applied to $S_{t} f_{1} \otimes \cdots \otimes S_{t} f_{n}$, using the differential form of the equation satisfied by $S_{t}$, Theorem 3.2(e). Taking the expectation and differentiating with respect to $t$, we get

$$
\begin{aligned}
\left.\frac{d}{d t}\right|_{t=0} P_{t}^{(n)} f_{1} \otimes \cdots \otimes f_{n} & =\left.\frac{d}{d t}\right|_{t=0} E\left[S_{t}^{\otimes n} f_{1} \otimes \cdots \otimes f_{n}\right] \\
& =A^{(n)} f_{1} \otimes \cdots \otimes f_{n}
\end{aligned}
$$

REMARK 4.4. In general, $m^{\otimes n}$ is not invariant under $P_{t}^{(n)}$.
5. Measure-preserving case. We say that the statistical solution $S_{t}$ is measure preserving if and only if $m S_{t}=m$ a.s. for all $t$ (i.e., $m$ is invariant for $S_{t}$ ). When $m(X)=\infty$, we use the natural extension of $S_{t}$ to $L^{1}(m)$ or to positive functions defined $m$-a.e.

Let us denote by $\mathscr{F}_{K}$ the set of functions of $\mathscr{F}$ which have compact support.
Proposition 5.1. $S_{t}$ is measure preserving if and only if $\int C(f, g) d m^{\otimes 2}$ vanishes for all $f, g$ in $\mathcal{F}_{K}$. Moreover, define $r_{t}$ on $L^{2}\left(\mathcal{F}_{t}\right)$ by $W_{s}^{k} \circ r_{t}=$ $W_{t-s}^{k}-W_{t}^{k}$. Then the adjoint of $S_{t}$ in $L^{2}(m)$ is $S_{t}^{*}=S_{t} \circ r_{t}$.

REMARK 5.2. (a) When $f \in \mathcal{F}_{K}, C(f, f) \in L^{1}\left(m^{\otimes 2}\right)$.
(b) In the Riemannian manifold case, the condition that $\int C(f, g) d m^{\otimes 2}$ vanishes for all $f, g$ in $\mathscr{F}_{K}$ is equivalent to assuming that $W_{t}$ is divergence free in the weak sense, that is, that for any $f \in \mathcal{F}_{K}, \int\left\langle W_{t}, \nabla f\right\rangle d m=0$. (It follows from the identity $E\left[\left(\int\left\langle W_{t}, \nabla f\right\rangle d m\right)^{2}\right]=t \int C(f, f) d m^{\otimes 2}$.)

LEMMA 5.3. Assume that $\int C(f, g) d m^{\otimes 2}$ vanishes for all $f, g$ in $\mathcal{F}_{K}$. Then, for every $h \in H, f, g$ in $\mathcal{F}$,

$$
\begin{equation*}
\int g D_{h} f d m=-\int f D_{h} g d m \tag{5.1}
\end{equation*}
$$

Proof. For every $h \in H,(g, f) \mapsto \int g D_{h} f d m$ is a continuous bilinear form on $\mathcal{F} \times \mathcal{F}$ since $\left\|D_{h} f\right\|_{L^{2}(m)}^{2} \leq \mathcal{E}(f, f)\|h\|_{L_{2}(m)}^{2}$.

Take $f, g$ in $\mathcal{F}_{K} \cap L^{\infty}(m)$. Then $f g \in \mathcal{F}_{K}$ (as the bounded functions of a Dirichlet space form an algebra) and, since $D_{e_{k}}$ is a derivation, $D_{e_{k}}(f g)=$ $g D_{e_{k}} f+f D_{e_{k}} g$. Using this property, we get

$$
\begin{aligned}
\sum_{k}\left(\int\left(g D_{e_{k}} f+f D_{e_{k}} g\right) d m\right)^{2} & =\sum_{k}\left(\int D_{e_{k}}(f g) d m\right)^{2} \\
& =\int C(f g, f g) d m^{\otimes 2}=0
\end{aligned}
$$

This implies that, for every $k, \int g D_{e_{k}} f d m=-\int f D_{e_{k}} g d m$. To conclude we observe that both members of (5.1) are continuous in $f$ and $g$ and that $\mathcal{F}_{K} \cap$ $L^{\infty}(m)$ is dense in $\mathcal{F}$ (since the Dirichlet form is regular; see [11], Section 1.1).

Proof of Proposition 5.1. Assume $\int C(f, g) d m^{\otimes 2}=0$ holds for every $f$ and $g$ in $\mathcal{F}_{K}$.

Let us remark that the expression of the $n$th chaos of $S_{t} f$ is given by the expression

$$
\begin{gather*}
J_{t}^{n} f=\int_{0 \leq s_{1} \leq s_{2} \leq \cdots \leq s_{n} \leq t} \sum_{k_{1}, \ldots, k_{n}} P_{s_{1}} D_{e_{k_{1}}} P_{s_{2}-s_{1}} D_{e_{k_{2}}} \cdots D_{e_{k_{n}}} P_{t-s_{n}} f  \tag{5.2}\\
\times d W_{s_{1}}^{k_{1}} \cdots d W_{s_{n}}^{k_{n}}
\end{gather*}
$$

From this expression, using Lemma 5.3 and the fact that $P_{t}$ is self-adjoint in $L^{2}(m)$, we get that, for $f$ and $g$ in $L^{2}(m)$,

$$
\begin{aligned}
& \int g J_{t}^{n} f d m \\
& (5.3)=\iint_{0 \leq s_{1} \cdots \leq s_{n} \leq t} f \sum_{k_{1}, \ldots, k_{n}}(-1)^{n} P_{t-s_{n}} D_{e_{k_{n}}} P_{s_{n}-s_{n-1}} \cdots P_{s_{2}-s_{1}} D_{e_{k_{1}}} P_{s_{1}} g \\
& \times d W_{s_{1}}^{k_{1}} \cdots d W_{s_{n}}^{k_{n}} d m
\end{aligned}
$$

Making the change of variable $u_{n-i+1}=t-s_{i}$, we get that the adjoint of $J_{t}^{n}$ is given by

$$
\begin{gather*}
\left(J^{n}\right)_{t}^{*} g=\int_{0 \leq u_{1} \leq u_{2} \leq \cdots \leq u_{n} \leq t} \sum_{k_{1}, \ldots, k_{n}} P_{u_{1}} D_{e_{k_{1}}} P_{u_{2}-u_{1}} D_{e_{k_{2}}} \cdots D_{e_{k_{n}}} P_{t-u_{n}} g \\
\times d W_{u_{1}}^{k_{1}} \circ r_{t} \cdots d W_{u_{n}}^{k_{n}} \circ r_{t} \tag{5.4}
\end{gather*}
$$

From this it is easy to see that $S_{t}^{*} g=\left(S_{t} \circ r_{t}\right) g$ (as they have the same chaos expansion).

Notice that $S_{t}^{*} 1=1$. A priori the constant functions are not in $L^{2}(m)$, but there exists an increasing sequence in $L^{2}(m), g_{n}$ such that $g_{n}$ converges toward 1 . For any nonnegative function $f \in L^{2}(m)$,

$$
\begin{equation*}
\int S_{t} f g_{n} d m=\int f S_{t}^{*} g_{n} d m \tag{5.5}
\end{equation*}
$$

This equation implies, taking the limit as $n$ goes to $\infty$, that

$$
\begin{equation*}
m S_{t}(f)=\int f S_{t}^{*} \mathbb{1} d m=m(f) \tag{5.6}
\end{equation*}
$$

and we get that $m S_{t}=m$ a.s., which ends the first part of the proof.
Conversely, it follows from Proposition 4.2 that, for all $f, g$ in $\mathscr{D}(A)$,

$$
S_{t}^{\otimes 2} f \otimes g-S_{t} f \otimes g-f \otimes S_{t} g+f \otimes g-\int_{0}^{t} S_{s}^{\otimes 2} C(f, g) d s
$$

is a square integrable martingale. This result extends to $f, g$ in $\mathcal{F}$. Taking $f, g$ in $\mathscr{F}_{K}$, integrating with respect to $m^{\otimes 2}$ and taking expectation, we get that $\int C(f, g) d m^{\otimes 2}$ vanishes.

REMARK 5.4. When $S_{t}$ is measure preserving, $P_{t}^{(n)}$ is self-adjoint in $L^{2}\left(m^{\otimes n}\right)$ and in particular $m^{\otimes n}$ is invariant under $P_{t}^{(n)}$. The associated local Dirichlet form $\mathcal{E}^{(2)}$ is such that

$$
\begin{aligned}
\mathcal{E}^{(2)}(f \otimes g, f \otimes g)= & \mathcal{E}(f, f)\|g\|_{L^{2}(m)}^{2}+\mathcal{E}(g, g)\|f\|_{L^{2}(m)}^{2} \\
& +2 \int C(f, g) f \otimes g d m^{\otimes 2}
\end{aligned}
$$

for any $(f, g) \in \mathcal{F}^{2}$ and a similar expression can be given for $\mathcal{E}^{(n)}$.
6. Existence of a flow of maps. Let $\left(S_{t}\right)_{t \geq 0}$ denote the statistical solution.

DEFINITION 6.1. We say that $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps if and only if there exists a family of measurable mappings $\left(\varphi_{t}\right)_{t \geq 0}$ from $X \times \Omega$ in $X$ such that, for any $f \in L^{2}(m)$ and any positive $t, S_{t} f=f \circ \bar{\varphi}_{t}$.

Note that if $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps, $P_{x, w}$ is the Dirac measure on the path $\left\{\varphi_{t}(x) ; t \geq 0\right\}$.

DEFINITION 6.2. We say that $\left(S_{t}\right)_{t \geq 0}$ is a coalescent flow of maps if and only if $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps and, for every $(x, y) \in X^{2}$, with positive probability there exists $T$ such that $\varphi_{t}(x)=\varphi_{t}(y)$ for all $t \geq T$.

Let $\left(\left(X_{t}, Y_{t}\right)\right)_{t \geq 0}$ designate the two-point motion associated with the statistical solution.

DEFINITION 6.3. We say that $\left(S_{t}\right)_{t \geq 0}$ is diffusive without hitting if and only if $\left(S_{t}\right)_{t \geq 0}$ is not a flow of maps and starting from $(x, x)$, for all positive $t, X_{t} \neq Y_{t}$.

DEFINITION 6.4. We say that $\left(S_{t}\right)_{t \geq 0}$ is diffusive with hitting if and only if $\left(S_{t}\right)_{t \geq 0}$ is not a flow of maps and $\left(X_{t}^{-}, Y_{t}\right)_{t \geq 0}$ hits the diagonal with positive probability.

In this section, we give conditions under which the statistical solution is a flow of maps or not.

LEMMA 6.5. $\quad\left(S_{t}\right)_{t \geq 0}$ is a flow of maps if and only if, for any $f \in L^{2}(m)$ and any positive $t, E\left[\left(S_{t} f\right)^{2}\right]=P_{t} f^{2}$.

Proof. It is clear that there exist Markovian kernels on $X, s_{t}(x, \omega, d y)$, such that $S_{t} f(x)=\int f(y) s_{t}(x, \omega, d y)$. Also, $s_{t}(x, \omega, d y)$ is the law of $X_{t}\left(\omega^{\prime}\right)$ under $P_{x, \omega}\left(d \omega^{\prime}\right)$. As $m \otimes P$-a.e,

$$
\begin{equation*}
\left(S_{t} f^{2}\right)(x)-\left(S_{t} f\right)^{2}(x)=\int\left(f(y)-\int f(z) s_{t}(x, \omega, d z)\right)^{2} s_{t}(x, \omega, d y) \tag{6.1}
\end{equation*}
$$

if $E\left[\left(S_{t} f\right)^{2}\right]=P_{t} f^{2}, \quad \int\left(f(y)-\int f(z) s_{t}(x, \omega, d z)\right)^{2} s_{t}(x, \omega, d y)=0 \quad$ and $s_{t}(x, \omega, d z)$ is a Dirac measure $\delta_{\varphi_{t}(x, \omega)}$, where $\varphi_{t}(x, \omega)$ is defined $m \otimes P$-a.e.

Let $h \in L^{1}(m)$ be a positive function such that $\int h d m=1$. For any positive $t$, let $\mu_{t}$ be a probability on the Borel sets of $X \times X$ such that, for any $(f, g) \in$ $L^{2}(m) \times L^{2}(m), \mu_{t}(f \otimes g)=\int E\left[S_{t} f S_{t} g\right] h d m$.

REMARK 6.6. $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps if and only if, for all positive $t$, $\mu_{t}(\Delta)=1$, where $\Delta=\{(x, x) ; x \in X\}$.

Proof. If $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps, there exists $\varphi_{t}$ such that $S_{t} f=f \circ \varphi_{t}$. If $A$ and $B$ are disjoint Borel sets of finite measure,

$$
\mu_{t}(A \times B)=\int E\left[\mathbb{1}_{A}\left(\varphi_{t}(x)\right) \mathbb{1}_{B}\left(\varphi_{t}(x)\right)\right] h(x) d m(x)=0
$$

This implies that $\mu_{t}(X \times X-\Delta)=0$ and as $\mu_{t}$ is a probability that $\mu_{t}(\Delta)=1$.
If $\mu_{t}(\Delta)=1$, for $f \in L^{2}(m), \mu_{t}\left(f^{2} \otimes \mathbb{1}-2 f \otimes f+1 \otimes f^{2}\right)=0$. This implies that

$$
\begin{equation*}
\int_{X} P_{t} f^{2} h d m=\int_{X} E\left[\left(S_{t} f\right)^{2}\right] h d m \tag{6.2}
\end{equation*}
$$

and that $E\left[\left(S_{t} f\right)^{2}\right]=P_{t} f^{2}$. Hence $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps.
Recall that we denoted by $P_{(\cdot, \cdot)}^{(2)}$ the law of the two-point motion $\left(\left(X_{t}, Y_{t}\right)\right)_{t \geq 0}$.

PROPOSITION 6.7. $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps if, for any positive $r$ and any positive $t$,

$$
\lim _{y \rightarrow x} P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r\right]=0, \quad m(d x) \text {-a.e. }
$$

Proof. For $\varepsilon>0$, let $v_{\varepsilon}$ be the measure on $X \times X$ such that, for any $(f, g) \in$ $L^{2}(m) \times L^{2}(m), v_{\varepsilon}(f \otimes g)=\int f P_{\varepsilon} g h d m$. For any $(f, g) \in L^{2}(m) \times L^{2}(m)$,

$$
\begin{equation*}
v_{\varepsilon} P_{t}^{(2)}(f \otimes g)=\int E\left[S_{t} f P_{\varepsilon} S_{t} g\right] h d m \tag{6.3}
\end{equation*}
$$

As $P_{\varepsilon} S_{t} g$ converges in $L^{2}(m \otimes P)$ toward $S_{t} g$ (see Remark 3.3),

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} v_{\varepsilon} P_{t}^{(2)}(f \otimes g)=\int E\left[S_{t} f S_{t} g\right] h d m \tag{6.4}
\end{equation*}
$$

Therefore, the family of measure $\left(v_{\varepsilon} P_{t}^{(2)}\right)_{\varepsilon>0}$ converges weakly toward $\mu_{t}$ as $\varepsilon$ goes to 0 .

Assume that, for any positive $r$ and any $t, \lim _{y \rightarrow x} P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r\right]=0$. Let $A$ and $B$ be two disjoint Borel sets such that $d(A, B) \geq r$. Then

$$
v_{\varepsilon} P_{t}^{(2)}(A \times B)=\int_{X} f_{\varepsilon}(x) h(x) d m(x)
$$

with

$$
f_{\varepsilon}(x)=\int P_{(x, y)}^{(2)}\left[X_{t} \in A \text { and } Y_{t} \in B\right] p_{\varepsilon}(x, d y)
$$

where $p_{\varepsilon}(x, d y)$ is the kernel given by $P_{\varepsilon}$.
As $d(A, B) \geq r$,

$$
f_{\varepsilon}(x) \leq \int P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r\right] p_{\varepsilon}(x, d y)
$$

For any positive $\beta$, for $m$-almost every $x$, there exists $\alpha(x)$ such that $d(x, y) \leq$ $\alpha(x)$ implies that $P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r\right] \leq \beta$. Note that

$$
f_{\varepsilon}(x) \leq \int_{\{d(x, y)>\alpha(x)\}} p_{\varepsilon}(x, d y)+\int_{\{d(x, y) \leq \alpha(x)\}} P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r\right] p_{\varepsilon}(x, d y)
$$

It is clear that

$$
\lim _{\varepsilon \rightarrow 0} \int_{\{d(x, y)>\alpha(x)\}} p_{\varepsilon}(x, d y)=0, \quad m(d x) \text {-a.e. }
$$

Hence, $\lim \sup f_{\varepsilon}(x) \leq \beta m(d x)$-a.e. and this holds for any positive $\beta$. Therefore, $\lim _{\varepsilon \rightarrow 0} f_{\varepsilon}(x)=0 m(d x)$-a.e. and, by dominated convergence $\left(\left|f_{\varepsilon}(x)\right| \leq 1\right)$,

$$
\lim _{\varepsilon \rightarrow 0} v_{\varepsilon} P_{t}^{(2)}(A \times B)=0
$$

This implies that $\mu_{t}(X \times X-\Delta)=0$ and that $\left(S_{t}\right)_{t \geq 0}$ is a flow of maps.

PROPOSITION 6.8. If there exist a positive $t$, a positive $r$ and $p \in] 0,1]$ such that, for $m^{\otimes 2}$-almost every $(x, y), P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right)>r\right] \geq p$, then $\left(S_{t}\right)_{t \geq 0}$ is not a flow of maps.

Proof. Suppose there exist a positive $t$, a positive $r$ and $p \in] 0,1]$ such that for $m^{\otimes 2}$-almost every $(x, y), P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right)>r\right] \geq p$.

Let $\left(B_{i}\right)_{i \in \mathbb{N}}$ be a partition of $X$ such that the diameter of $B_{i}$ is lower than $r$.
Let us suppose that $\mu_{t}(\Delta)=1$ [or that $\left(S_{t}\right)_{t}$ is a flow of maps]. Then we have $\sum_{i} \mu_{t}\left(B_{i} \times B_{i}\right)=1$, and for any positive $\alpha$, there exists $N$ such that

$$
\sum_{i=1}^{N} \mu_{t}\left(B_{i} \times B_{i}\right) \geq 1-\alpha
$$

Since $v_{\varepsilon} P_{t}^{(2)}$ converges weakly toward $\mu_{t}$,

$$
\begin{aligned}
\sum_{i=1}^{N} & \mu_{t}\left(B_{i} \times B_{i}\right) \\
& =\lim _{\varepsilon \rightarrow 0} v_{\varepsilon} P_{t}^{(2)}\left(B_{i} \times B_{i}\right) \\
& =\lim _{\varepsilon \rightarrow 0} \sum_{i=1}^{N} \int_{X \times X} P_{(x, y)}^{(2)}\left[\left(X_{t}, Y_{t}\right) \in B_{i} \times B_{i}\right] p_{\varepsilon}(x, d y) h(x) d m(x) \\
& \leq \lim _{\varepsilon \rightarrow 0} \sum_{i=1}^{N} \int_{X \times X} P_{(x, y)}^{(2)}\left[X_{t} \in B_{i} ; d\left(X_{t}, Y_{t}\right) \leq r\right] p_{\varepsilon}(x, d y) h(x) d m(x) \\
& \leq \lim _{\varepsilon \rightarrow 0} \int_{X \times X} P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \leq r\right] p_{\varepsilon}(x, d y) h(x) d m(x) \leq 1-p
\end{aligned}
$$

Choosing $\alpha<p$, we get a contradiction. Hence $\mu_{t}(\Delta)<1$ and $\left(S_{t}\right)_{t \geq 0}$ is not a flow of maps.
7. A one-dimensional example. Let $X=\mathbb{R}$, let $P_{t}$ be the semigroup of the Brownian motion on $\mathbb{R}$ and let the covariance function $C(x, y)=\operatorname{sgn}(x) \operatorname{sgn}(y)$ [where $\operatorname{sgn}(x)$ denotes the sign of $x$ with the convention $\operatorname{sgn}(0)=1$ ]. Here, we have $W_{t}(x)=\operatorname{sgn}(x) W_{t}$, where $W_{t}$ is a Brownian motion starting from 0 . Set $L_{t}^{x}=\sup _{s \leq t}\left\{-\operatorname{sgn}(x)\left(x+W_{s}\right)\right\} \vee 0$ and $R_{t}^{x}=x+W_{t}+\operatorname{sgn}(x) L_{t}^{x}$ (it is a Brownian motion starting from $x$, reflected at 0 ).

Proposition 7.1. The statistical solution $S_{t}$ can be written as

$$
\begin{equation*}
S_{t} f(x)=f\left(R_{t}^{x}\right) \mathbb{1}_{L_{t}^{x}=0}+\frac{1}{2}\left[f\left(R_{t}^{x}\right)+f\left(-R_{t}^{x}\right)\right] \mathbb{1}_{L_{t}^{x}>0} . \tag{7.1}
\end{equation*}
$$

Proof. On an extension of the probability space, it is possible to build a Brownian motion starting from $x, X_{t}$, such that $W_{t}=\int_{0}^{t} \operatorname{sgn}\left(X_{s}\right) d X_{s}$ [then $X_{t}$ is a weak solution of the $\left.\operatorname{SDE} d X_{t}=\operatorname{sgn}\left(X_{t}\right) d W_{t}\right]$. Then $S_{t} f(x)=E\left[f\left(X_{t}\right) \mid \mathcal{F}^{B}\right]$, with $\mathcal{F}^{B}=\sigma\left(W_{u} ; u \geq 0\right)$. Let us remark that $L_{t}^{x}$ is the local time of $X$ at 0 and that $R_{t}^{x}=\operatorname{sgn}(x)\left|X_{t}\right|$. Set $T=\inf \left\{t ; L_{t}^{x}>0\right\}=\inf \left\{t ; X_{t}=0\right\}$. Formula (7.1) follows simply from the fact that

$$
E\left[f\left(X_{t}\right) \mathbb{1}_{t \geq T}| | X_{t} \mid\right]=\frac{1}{2}\left(f\left(X_{t}\right)+f\left(-X_{t}\right)\right) \mathbb{1}_{t \geq T}
$$

8. The Lipschitz case. Assume $X$ is a Riemannian manifold with injectivity radius $\rho>0$. Let $P_{t}$ be the semigroup of a symmetric diffusion on $X$ with generator $A$. Let $C$ be a covariance inducing the metric [i.e., with equality in (1.7)].

We will say that $C$ is Lipschitz if and only if there exist a positive constant $k$ and $0<\varepsilon<\rho$ such that, for any $(x, y) \in X^{2}$, with $d(x, y)<\varepsilon$,

$$
\begin{equation*}
A^{(2)} d^{2}(x, y) \leq k d^{2}(x, y) \tag{8.1}
\end{equation*}
$$

REMARK. (a) $d^{2}(x, y)$ is smooth on $\left\{(x, y) \in X^{2}, d(x, y)<\rho\right\}$ since $\rho$ is the injectivity radius.
(b) On $\mathbb{R}^{d}$, the condition (8.1) will be checked as soon as

$$
\begin{gather*}
A=\frac{1}{2} \sum_{1 \leq i, j \leq d}^{d} C^{i j}(x, x) \partial_{i} \partial_{j}+\sum_{i} b^{i}(x) \partial_{i} \\
\sum_{i=1}^{d}\left(C^{i i}(x, x)+C^{i i}(y, y)-2 C^{i i}(x, y)\right) \leq \frac{k}{2} d(x, y)^{2} \tag{8.2}
\end{gather*}
$$

and $b^{i}$ is a Lipschitz function for all $i$.
Equation (8.2) is satisfied when $C$ is $C^{2}$ or when $C=\sum_{\alpha=1}^{n} X_{\alpha} \otimes X_{\alpha}$, where $X_{\alpha}$ are Lipschitz vector fields. In the latest case, the flow of maps can be constructed by the usual fixed point method for solutions of SDEs based on Gronwall's lemma.

Let $\left(X_{t}, Y_{t}\right)$ be the two-point motion associated with the statistical solution. Set $\tau=\inf \left\{t, d\left(X_{t}, Y_{t}\right) \geq \varepsilon\right\}$ and $H_{t}=d^{2}\left(X_{t \wedge \tau}, Y_{t \wedge \tau}\right)$.

LEMMA 8.1. $\quad E_{(x, y)}^{(2)}\left(H_{t}\right) \leq e^{k t} d^{2}(x, y)$.
Proof. By Itô's formula,

$$
H_{t}-H_{0}=M_{t}+\int_{0}^{t \wedge \tau} A^{(2)} d^{2}\left(X_{s}, Y_{s}\right) d s
$$

where $M_{t}$ is a martingale. Hence

$$
\begin{aligned}
H_{t}-H_{0} & \leq M_{t}+\int_{0}^{t \wedge \tau} k d^{2}\left(X_{s}, Y_{s}\right) d s \\
& \leq M_{t}+\int_{0}^{t} k H_{s} d s
\end{aligned}
$$

This implies that $E_{(x, y)}^{(2)}\left(H_{t}\right)-d^{2}(x, y) \leq k \int_{0}^{t} E_{(x, y)}^{(2)}\left(H_{s}\right) d s$. Hence the lemma.

THEOREM 8.2. Assume (8.1) is satisfied. Then the statistical solution associated with $P_{t}$ and $C$ is a flow of maps.

Proof. Indeed, for any $r<\varepsilon$,

$$
\begin{aligned}
P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r\right] & \leq P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right) \geq r \text { or } t \geq \tau\right] \leq \frac{1}{r^{2}} E_{(x, y)}^{(2)}\left(H_{t}\right) \\
& \leq \frac{e^{k t}}{r^{2}} d(x, y)^{2}
\end{aligned}
$$

which goes to 0 as $d(x, y)$ goes to 0 . We conclude using Theorem 6.7.

## 9. Isotropic statistical solution on $S^{d}$.

9.1. Isotropic covariance function on $S^{d}$. On $S^{d}$ with $d \geq 2$, isotropic covariance functions are given by the formula (see [25])

$$
\begin{equation*}
C((x, u),(y, v))=\alpha(t)\langle u, v\rangle+\beta(t)\langle u, y\rangle\langle v, x\rangle, \tag{9.1}
\end{equation*}
$$

with $(x, y) \in S^{d} \times S^{d}, t=\langle x, y\rangle=\cos \varphi$ and $(u, v) \in T_{x} S^{d} \times T_{y} S^{d} . \alpha$ and $\beta$ are given by

$$
\begin{align*}
& \alpha(t)=\sum_{l=1}^{\infty} a_{l} \gamma_{l}(t)+\sum_{l=1}^{\infty} b_{l}\left(t \gamma_{l}(t)-\frac{1-t^{2}}{d-1} \gamma_{l}^{\prime}(t)\right),  \tag{9.2}\\
& \beta(t)=\sum_{l=1}^{\infty} a_{l} \gamma_{l}^{\prime}(t)+\sum_{l=1}^{\infty} b_{l}\left(-\gamma_{l}(t)-\frac{t}{d-1} \gamma_{l}^{\prime}(t)\right), \tag{9.3}
\end{align*}
$$

where $\gamma_{l}(t)=C_{l-1}^{(d+1) / 2}(t) / C_{l-1}^{(d+1) / 2}(1), C_{l}^{p}$ is a Gegenbauer polynomial, and $a_{l}$ and $b_{l}$ are nonnegative such that $\sum_{l} a_{l}<\infty$ and $\sum_{l} b_{l}<\infty$. Using the integral form of the Gegenbauer polynomials (see [27], page 496),

$$
\begin{equation*}
\gamma_{l}(\cos \varphi)=\int_{0}^{\pi}[z(\varphi, \theta)]^{l-1} \sin ^{d} \theta \frac{d \theta}{c_{d}} \tag{9.4}
\end{equation*}
$$

with $c_{d}=\int_{0}^{\pi} \sin ^{d} \theta d \theta$ and $z(\varphi, \theta)=\cos \varphi-i \sin \varphi \cos \theta$.
In [12], it is proved that the spectrum of the Laplacian $\Delta$ acting on the $L^{2}$-vector fields is $\{-l(l+d-1), l \geq 1\} \cup\{-(l+1)(l+d-2), l \geq 1\}$. Let $\mathscr{g}_{l}$ and $\mathscr{D}_{l}$ be respectively the eigenspaces corresponding to the eigenvalues $-l(l+d-1)$ and $-(l+1)(l+d-2) . g_{l}$ is constituted of gradient vector fields and $\mathscr{D}_{l}$ of divergence-free vector fields. These spaces can be isometrically identified with the spaces $\mathscr{H}_{d+1, l}$ and $\mathcal{F}_{d+1, l}$ used in [25] and can be used as carrier spaces of the irreducible representations of $S O(d+1), T^{l}$ and $Q^{l}$.

Let $\left(\alpha_{M}^{l}\right)_{M}$ and $\left(\omega_{M}^{l}\right)_{M}$ be orthonormal bases of $\mathscr{g}_{l}$ and $\mathscr{D}_{l}$. Then, if $\left(z_{M, d}^{l}\right)_{l, M}$ and $\left(z_{M, \delta}^{l}\right)_{l, M}$ are independent families of independent normalized centered Gaussian variables,

$$
\begin{equation*}
W=\sum_{l \geq 1} \sqrt{\frac{d a_{l}}{\operatorname{dim} \mathscr{g}_{l}}} \sum_{M} z_{M, d}^{l} \alpha_{M}^{l}+\sum_{l \geq 1} \sqrt{\frac{d b_{l}}{\operatorname{dim} \mathscr{D}_{l}}} \sum_{M} z_{M, \delta}^{l} \omega_{M}^{l} \tag{9.5}
\end{equation*}
$$

is an isotropic Gaussian vector field of covariance $C$ given by (9.1), (9.2) and (9.3).

SkETCH OF PROOF. The covariance of $W$ is

$$
\sum_{l \geq 1} \frac{d a_{l}}{\operatorname{dim} \mathcal{g}_{l}} \sum_{M} \alpha_{M}^{l} \otimes \alpha_{M}^{l}+\sum_{l \geq 1} \frac{d b_{l}}{\operatorname{dim} \mathscr{D}_{l}} \sum_{M} \omega_{M}^{l} \otimes \omega_{M}^{l}
$$

Let us choose $\left(\alpha_{M}^{l}\right)_{M}$ such that $\alpha_{M}^{l}=c_{1}(l, d) \nabla \Xi_{M}^{l}$ (where $\left(\Xi_{M}^{l}\right)_{M}$ is the basis of $\mathscr{H}_{d+1, l}$ given in [25]). Then, using the fact that $\Xi_{M}^{l}(p)=0$ if $M \neq 0$, for $x=g_{1} p$ and $y=g_{2} p[$ with $p=(0, \ldots, 0,1)]$,

$$
\begin{aligned}
\sum_{M} \Xi_{M}^{l}(x) \Xi_{M}^{l}(y) & =\sum_{M, N, K} T_{M N}^{l}\left(g_{1}\right) T_{M K}^{l}\left(g_{2}\right) \Xi_{N}^{l}(p) \Xi_{K}^{l}(p) \\
& =T_{00}^{l}\left(g_{2}^{-1} g_{1}\right)\left(\Xi_{0}^{l}(p)\right)^{2}
\end{aligned}
$$

In [27, 25], $T_{00}^{l}(g)$ is computed and it is easy from this to give the covariance of the gradient part of $W$. We can calculate the covariance of the divergencefree part similarly. We choose the orthonormal basis $\left(\omega_{M}^{l}\right)_{M}$ of $\mathscr{D}_{l}$ such that, for $M \notin\{1, \ldots, d\}, \omega_{M}^{l}(p)=0$ and such that, for $1 \leq i \leq d, \omega_{i}^{l}(p)=c_{2}(l, d) e_{i}$ (this basis corresponds to the basis of $\mathcal{F}_{d+1, l}$ given in [25]). Then one has, for $x=g p$ and $g \in S O(d)$,

$$
\begin{equation*}
\omega_{M}^{l}(x)=\sum_{i=1}^{d} Q_{M i}^{l}(g) g\left(\omega_{i}^{l}(p)\right)=c_{2}(l, d) Q_{M i}^{l}(g) g\left(e^{i}\right) \tag{9.6}
\end{equation*}
$$

Then, for every $(x, u)$ and $(y, v)$ in $T S^{d}$,

$$
\begin{align*}
& \sum_{M}\left\langle\omega_{M}^{l}(x), u\right\rangle\left\langle\omega_{M}^{l}(y), v\right\rangle \\
& \quad=\left(c_{2}(l, d)\right)^{2} \sum_{M} Q_{M i}^{l}\left(g_{1}\right) Q_{M j}^{l}\left(g_{2}\right)\left\langle g_{1}\left(e^{i}\right), u\right\rangle\left\langle g_{2}\left(e^{j}\right), v\right\rangle  \tag{9.7}\\
& \quad=\left(c_{2}(l, d)\right)^{2} Q_{j i}^{l}(g)\left\langle g_{1}\left(e^{i}\right), u\right\rangle\left\langle g_{2}\left(e^{j}\right), v\right\rangle \tag{9.8}
\end{align*}
$$

with $g=g_{2}^{-1} g_{1}$. In [25], the matrix elements $Q_{j i}^{l}(g)$ are calculated and it is easy from this to give the covariance of the divergence-free part of $W$.

Let us now introduce Sobolev spaces and related covariances.
Let $H^{2, s}$ be the Sobolev space obtained by completion of the smooth vector fields with respect to the norm $\left\langle\left(-\Delta+m^{2}\right)^{s} V, V\right\rangle_{2}$ [with $\langle V, V\rangle_{2}=$ $\int\|V(x)\|^{2} d x$ ], where $m$ is positive. Note that the definition of $H^{2, s}$ does not depend on $m$.

Let $a$ and $b$ be nonnegative reals. Take $a_{l}=\frac{a}{(l-1)^{\alpha+1}}$ and $b_{l}=\frac{b}{(l-1)^{\alpha+1}}$ for $l \geq 1$ and $a_{1}=b_{1}=0$. For $\alpha>0$, set $G(\varphi)=\sum_{l \geq 2} \frac{1}{(l-1)^{\alpha+1}} \gamma_{l}(\cos \varphi)$. The function $G$ is well defined on $[0, \pi]$ as $\left|\gamma_{l}\right| \leq 1$.

Let $F_{d}$ and $F_{\delta}$ be real functions such that, for all $l \geq 2$,

$$
\begin{array}{r}
(l-1)^{\alpha+1} \operatorname{dim} \mathcal{g}_{l} \cdot F_{d}(-l(l+d-1))=d \\
(l-1)^{\alpha+1} \operatorname{dim} \mathscr{D}_{l} \cdot F_{\delta}(-(l+1)(l+d-2))=d \tag{9.10}
\end{array}
$$

and $F_{d}(-d)=F_{\delta}(-2(d-1))=0$. Note that, when $d=2, F_{d}=F_{\delta}$.
Let $\Pi$ be the orthonormal projection on the space of the $L^{2}$-gradient vector fields.

Proposition 9.1. The covariance function defined by the sequences $\left(a_{l}\right)$ and $\left(b_{l}\right)$ is given by (9.1) with the functions

$$
\begin{align*}
& \alpha(\cos \varphi)=a G(\varphi)+b\left(\cos \varphi G(\varphi)+\frac{\sin \varphi}{d-1} G^{\prime}(\varphi)\right),  \tag{9.11}\\
& \beta(\cos \varphi)=-\frac{a}{\sin \varphi} G^{\prime}(\varphi)+b\left(-G(\varphi)+\frac{\cos \varphi}{(d-1) \sin \varphi} G^{\prime}(\varphi)\right) . \tag{9.12}
\end{align*}
$$

When $a$ and $b$ are positive, the associated self-reproducing space is $H^{2,(\alpha+d) / 2}$ equipped with a different (but equivalent) norm, namely

$$
\|V\|_{H}^{2}=\frac{1}{a}\|\Pi V\|_{d}^{2}+\frac{1}{b}\|(I-\Pi) V\|_{\delta}^{2}
$$

where $\|V\|_{d}^{2}=\left\langle F_{d}(\Delta)^{-1} V, V\right\rangle_{2}$ and $\|V\|_{\delta}^{2}=\left\langle F_{\delta}(\Delta)^{-1} V, V\right\rangle_{2}$.
Proof. It is not difficult to see that the norm $\|\cdot\|_{H}$ given in the proposition is the norm on the self-reproducing space associated with $C$.

Now since (see [12])

$$
\begin{aligned}
\operatorname{dim} g_{l} & =\frac{(d+l-3)!}{(d-1)!(l-1)!}(d+2 l-3)(d+1) \\
\operatorname{dim} \mathscr{D}_{l} & =\frac{(d+l-3)!}{(d-1)!(l-1)!}(d+2 l-3) \frac{d(d+1)}{2}
\end{aligned}
$$

for $\lambda \rightarrow \infty, \lambda^{(\alpha+d) / 2} F_{d}(\lambda)=O(1)$ and $\lambda^{(\alpha+d) / 2} F_{\delta}(\lambda)=O(1)$. This implies that $\|\cdot\|_{H}$ and the norm used to define $H^{2,(\alpha+d) / 2}$ are equivalent (when $a$ and $b$ are positive). We get that the self-reproducing space associated with $C$ is $H^{2,(\alpha+d) / 2}$.

REMARK 9.2. If $a$ or $b$ vanishes, the self-reproducing space is $H^{2,(\alpha+d) / 2}$ restricted to divergence-free vector fields or gradient vector fields.
9.2. Phase transitions for the Sobolev statistical solution. Let $P_{t}$ be the semigroup of the Brownian motion of variance $(a+b) G(0)$ and $S_{t}$ be the statistical solution associated with $P_{t}$ and $C$.

Let $\left(X_{t}, Y_{t}\right)$ be the two-point motion. Let $\psi_{t}=d\left(X_{t}, Y_{t}\right)$. Since $h(x, y)=$ $d^{2}(x, y)$ is a $C^{2}$-function, $h$ belongs to $\mathscr{D}\left(A^{(2)}\right)$ and since $X_{t}$ and $Y_{t}$ are solutions of an SDE like (3.3), $\psi_{t}^{2}$ is a diffusion on $\left[0, \pi^{2}\right]$ and is a solution of an SDE; $\psi_{t}$ is also a diffusion on $[0, \pi]$ [note that $d(x, y)$ a priori does not belong to $\left.\mathscr{D}\left(A^{(2)}\right)\right]$. This diffusion is eventually reflected (or absorbed) in 0 and $\pi$. Its generator is $L=\sigma^{2}(\varphi) \frac{d^{2}}{d \varphi^{2}}+b(\varphi) \frac{d}{d \varphi}($ see [25]), with

$$
\begin{align*}
\sigma^{2}(\varphi) & =\alpha(1)-\alpha(\cos \varphi) \cos \varphi+\beta(\cos \varphi) \sin ^{2} \varphi,  \tag{9.13}\\
b(\varphi) & =\frac{(d-1)}{\sin \varphi}(\alpha(1) \cos \varphi-\alpha(\cos \varphi)) . \tag{9.14}
\end{align*}
$$

The generator of $\psi_{t}^{2}$ is $L^{\prime}=\widetilde{\sigma}^{2}(x) \frac{d^{2}}{d x^{2}}+\widetilde{b}(x) \frac{d}{d x}$, with

$$
\begin{align*}
\tilde{\sigma}^{2}(x) & =4 x \sigma^{2}(\sqrt{x}),  \tag{9.15}\\
\widetilde{b}(x) & =2 \sigma^{2}(\sqrt{x})+2 \sqrt{x} b(\sqrt{x}) . \tag{9.16}
\end{align*}
$$

Lemma 9.3. If $\alpha>2$, the statistical solution is a flow of maps.
Proof. We have $A^{(2)} d^{2}(x, y)=2 \sigma^{2}(d(x, y))+2 b(d(x, y)) d(x, y)$. When $\alpha>2$, then $G$ is $C^{2}$; this implies that $\alpha$ is $C^{2}$ and $\beta$ is continuous. Hence (8.1) can be checked.

Suppose $a+b>0$ and let $\eta=\frac{b}{a+b}$.
Theorem 9.4. For any $\alpha \in] 0,2[$, the following hold:
(a) For $d=2$ or 3 and $\eta<1-\frac{d}{\alpha^{2}}$, the statistical solution is a coalescent flow of maps.
(b) For $d=2$ or 3 and $1-\frac{d}{\alpha^{2}}<\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, the statistical solution is diffusive with hitting.
(c) For $d=2$ or 3 and $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$ or for $d \geq 4$, the statistical solution is diffusive without hitting.

REMARK. The same phase transition appears in the $\mathbb{R}^{d}$ case (see Theorem 10.1 below). It has been independently observed, in the context of the advection of a passive scalar, by Gawedzky and Vergassola [14].

Lemma 9.5. For $\alpha \in] 0,2[$, we have the following:
(a) $G$ is differentiable on $] 0, \pi[$;
(b) $\lim _{\varphi \rightarrow 0+} \frac{G(0)-G(\varphi)}{\varphi^{\alpha}}=\int_{0}^{\pi} \int_{0}^{\infty} \frac{\cos ^{2} \theta}{t^{2}+\cos ^{2} \theta} t^{\alpha-1} \sin ^{d} \theta \frac{d t d \theta}{\Gamma(\alpha+1) c_{d}}=K G(0)$;
(c) $\lim _{\varphi \rightarrow 0+} \frac{G^{\prime}(\varphi)}{\varphi^{\alpha-1}}=-\alpha \int_{0}^{\pi} \int_{0}^{\infty} \frac{\cos ^{2} \theta}{t^{2}+\cos ^{2} \theta} t^{\alpha-1} \sin ^{d} \theta \frac{d t d \theta}{\Gamma(\alpha+1) c_{d}}=-\alpha K G(0)$.

The proof of Lemma 9.5 is in Appendix A. From this lemma, we get, as $\varphi$ goes to 0 ,

$$
\begin{align*}
& \alpha(\cos \varphi)=(a+b) G(0)-\left(a+\left(1+\frac{\alpha}{d-1}\right) b\right) K G(0) \varphi^{\alpha}+o\left(\varphi^{\alpha}\right)  \tag{9.17}\\
& \beta(\cos \varphi)=\alpha\left(a-\frac{b}{d-1}\right) K G(0) \varphi^{\alpha-2}+o\left(\varphi^{\alpha-2}\right) \tag{9.18}
\end{align*}
$$

Hence,

$$
\begin{align*}
\sigma^{2}(\varphi) & =(a+b) K G(0)(\alpha+1-\alpha \eta) \varphi^{\alpha}(1+o(1))  \tag{9.19}\\
b(\varphi) & =(a+b) K G(0)(d-1+\alpha \eta) \varphi^{\alpha-1}(1+o(1)) \tag{9.20}
\end{align*}
$$

To prove Theorem 9.4, we need to study the two-point motion. Because of isotropy, it is enough to study the diffusion $\psi_{t}$. This diffusion satisfies an SDE until it exits $] 0, \pi[$.

Let $s$ be the scale function of the diffusion $\psi_{t}$,

$$
\left.s(x)=\int_{x_{0}}^{x} \exp \left[-\int_{x_{0}}^{y} \frac{b(\varphi)}{\sigma^{2}(\varphi)} d \varphi\right] d y \quad \text { with }\left(x_{0}, x\right) \in\right] 0, \pi\left[^{2}\right.
$$

Let $x \in\{0, \pi\}$ and $T_{x}=\inf \left\{t>0 ; \psi_{t}=x\right\}$. Using Breiman's terminology (see [4], pages 368-369), $x$ is an open boundary point if $T_{x}=\infty$ and is a closed boundary point if $T_{x}<\infty$. Note that $x$ is an open boundary point if $|s(x)|=\infty$.

First we are going to show that $\pi$ is an open boundary point. Then:

1. when $d=2$ or 3 and $\eta<1-\frac{d}{\alpha^{2}}$, we prove that 0 is an exit boundary point (this implies that the statistical solution is a coalescent flow of maps);
2. when $d=2$ or 3 and $1-\frac{d}{\alpha^{2}}<\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, we prove that 0 is an instantaneously reflecting regular boundary point (this implies that the statistical solution is diffusive with hitting);
3. when $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, we prove that 0 is an open entrance boundary point (this implies that the statistical solution is diffusive without hitting).

LEMMA 9.6. $\pi$ is an open boundary point.
Proof. It is easy to check that $s(\pi-)=\infty$ using the fact that $\alpha(1)+$ $\alpha(-1)>0$ :

$$
\begin{aligned}
\alpha(1)+\alpha(-1) & =(a+b) G(0)+(a-b) G(\pi) \\
& >(a+b) G(\pi)+(a-b) G(\pi) \geq 0
\end{aligned}
$$

Since $\pi$ is an open boundary point, we now study the behavior of $\psi_{t}$ at and near 0 .

LEMMA 9.7. If $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}, s(0+)=-\infty$ and if $\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, $s(0+)>-\infty$.

Proof. Let us note $\mu=\frac{d-1+\alpha \eta}{\alpha+1-\alpha \eta}$. Then we have that $\frac{b(\varphi)}{\sigma^{2}(\varphi)}=\frac{\mu}{\varphi}(1+o(1))$ and for any positive $\varepsilon$ there exist positive constants $C_{1}$ and $C_{2}$ such that, for $y \leq x_{0}$,

$$
\begin{equation*}
C_{1} y^{-\mu+\varepsilon} \leq \exp \left[-\int_{x_{0}}^{y} \frac{b(\varphi)}{\sigma^{2}(\varphi)} d \varphi\right] \leq C_{2} y^{-\mu-\varepsilon} \tag{9.21}
\end{equation*}
$$

From this, we see that $s(0+)=-\infty$ if $\mu>1$ (or if $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$ ) and $s(0+$ ) is finite if $\mu<1$ (or if $\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}$ ).

Lemmas 9.6 and 9.7 imply that (see [16], Theorem VI-3.1) if $\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, we have $T_{0}<\infty, T_{\pi}=\infty$ a.s. and if $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}, 0$ is an open boundary point and we have $\liminf \psi_{t}=0$ and $\limsup \psi_{t}=\pi$ a.s. $\left(\psi_{t}\right.$ is recurrent $)$.

REMARK 9.8. When $d \geq 4$ and $\alpha \in] 0,2\left[, \frac{1}{2}-\frac{(d-2)}{2 \alpha}<0\right.$. This implies that $\liminf \psi_{t}=0$ and $\limsup \psi_{t}=\pi$ a.s.

Since $\pi$ is an open boundary point, $\psi_{t} \in\left[0, \pi\right.$ [ for every positive $t$ and $\psi_{t}^{2}$ is a solution of the SDE

$$
\begin{equation*}
d \psi_{t}^{2}=\sqrt{2} \widetilde{\sigma}\left(\psi_{t}^{2}\right) d B_{t}+\widetilde{b}\left(\psi_{t}^{2}\right) d t \tag{9.22}
\end{equation*}
$$

Note that 0 is a solution of this $\operatorname{SDE}$ [since $\widetilde{\sigma}(0)=\widetilde{b}(0)=0$ ]. The solutions of this SDE might be not unique.

Let $m(d x)$ be the speed measure of the diffusion

$$
m(d x)=\mathbb{1}_{] 0, \pi[ }(x) \exp \left[\int_{x_{0}}^{x} \frac{b(\varphi)}{\sigma^{2}(\varphi)} d \varphi\right] \frac{d x}{\sigma^{2}(x)}+m(\{0\}) \delta_{0}=g(x) d x+m(\{0\}) \delta_{0}
$$

with $\left.x_{0} \in\right] 0, \pi[$.
LEMMA 9.9. If $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}, 0$ is an entrance open boundary point.
Proof. When $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}, 0$ is an open boundary point. From [4, Proposition 16.45], 0 is an entrance boundary point if and only if $\int_{0+}|s(x)| m(d x)<\infty$. For any positive $\varepsilon$, there exists a positive constant $D$ such that, for any $x \in] 0, x_{0}[$,

$$
|s(x) g(x)| \leq D x^{(e / c-\alpha-\varepsilon) \wedge 0} x^{-e / c-\varepsilon+1} \leq D x^{1-\alpha-2 \varepsilon}
$$

This shows that $\int_{0+} s(x) m(d x)<\infty$ (choose $\varepsilon$ such that $\left.2 \varepsilon \leq 2-\alpha\right)$.

This lemma implies that, when $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, there exist a positive $t$, a positive $\alpha$ and $p \in] 0,1[$ such that, for any $x \in] 0, \pi\left[, P_{x}\left[\psi_{t}>\alpha\right]>p\right.$. Proposition 6.8 implies that $S_{t}$ is not a flow of maps and since 0 is open, $S_{t}$ is diffusive without hitting.

Now let $d \in\{2,3\}$ (when $d \geq 4$ we always have $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$ ).
Lemma 9.10. If $\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}, 0$ is a closed boundary point.
Proof. From Proposition 16.43 in [4], page 366, $T_{0}$ is finite or the boundary point 0 is closed if, and only if, for any $b \in] 0, \pi\left[, \int_{0}^{b}|s(x)-s(0)| m(d x)\right.$ is finite.

We have

$$
|s(x)-s(0)| g(x) \sim \int_{0}^{x} \exp \left[-\int_{x_{0}}^{y} \frac{b(\varphi)}{\sigma^{2}(\varphi)} d \varphi\right] \frac{1}{\sigma^{2}(x)} \exp \left[\int_{x_{0}}^{y} \frac{b(\varphi)}{\sigma^{2}(\varphi)} d \varphi\right] d y
$$

Hence $|s(x)-s(0)| g(x)=O\left(x^{1-\alpha}\right)$. This implies that $\int_{0}^{b}|s(x)-s(0)| m(d x)$ is finite. This proves that $T_{0}$ is finite a.s.

Lemma 9.11. If $\eta<1-\frac{d}{\alpha^{2}}, 0$ is an exit boundary point.
Proof. In [4], 0 is an exit boundary point if and only if $m(] 0, x[)=\infty$ for all $x \in] 0, \pi\left[\right.$. This is the case if $\mu-\alpha<-1$ (or if $\eta<1-\frac{d}{\alpha^{2}}$ ). Note that, for $d=2$ or 3 and $\alpha \in] 0,2\left[, 1-\frac{d}{\alpha^{2}}<\frac{1}{2}-\frac{(d-2)}{2 \alpha}\right.$.

Lemma 9.11 implies that, when $\eta<1-\frac{d}{\alpha^{2}}$, the diffusion $\psi_{t}$ is absorbed at 0 , and, for any positive $r$,

$$
\lim _{d(x, y) \rightarrow 0} P_{(x, y)}^{(2)}\left[d\left(X_{t}, Y_{t}\right)>r\right]=\lim _{\varphi \rightarrow 0} P_{\varphi}\left[\psi_{t}>r\right]=0 .
$$

Now, applying Proposition 6.7, we prove that the statistical solution is a flow of maps and this is a coalescent flow of maps (since 0 is an exit boundary point).

Lemma 9.12. If $\eta \in] 1-\frac{d}{\alpha^{2}}, \frac{1}{2}-\frac{(d-2)}{2 \alpha}[, 0$ is a regular boundary point.
Proof. In [4], we see that 0 is regular if $m(] 0, x[)<\infty$ for all $x \in] 0, \pi[$, which is the case when $\eta \in] 1-\frac{d}{\alpha^{2}}, \frac{1}{2}-\frac{(d-2)}{2 \alpha}[$.

When $\eta \in] 1-\frac{d}{\alpha^{2}}, \frac{1}{2}-\frac{(d-2)}{2 \alpha}$ [, the two-point motion hits the diagonal. However, there is no uniqueness of the solution of the SDE satisfied by $\psi_{t}$ since 0 might be absorbing or (slowly or instantaneously) reflecting. To finish the proof of Theorem 9.4 , we prove that 0 is instantaneously reflecting.

To prove this, for $\varepsilon \in] 0,1\left[\right.$, let us introduce the covariance $C_{\varepsilon}=(1-\varepsilon)^{2} C$ [then, if $W_{t}$ is the cylindrical Brownian motion associated with $C,(1-\varepsilon) W_{t}$ is
the cylindrical Brownian motion associated with $C_{\varepsilon}$ ], and let $S_{t}^{\varepsilon}$ be the statistical solution associated with $P_{t}$ and $C_{\varepsilon}$.

For $f \in L^{2}(d x), S_{t}^{\varepsilon} f=\sum_{n \geq 0} J_{t}^{n, \varepsilon} f$, where $J_{t}^{n, \varepsilon} f$ is the $n$th chaos in the chaos expansion of $S_{t}^{\varepsilon} f$. [Note that $S_{t}^{\varepsilon}=Q_{\log (1-\varepsilon)} S_{t}$, where $Q_{\alpha}$ is the OrnsteinUhlenbeck operator on the Wiener space (used in Malliavin calculus; see [24]).] It is easy to see that $J_{t}^{n, \varepsilon} f=(1-\varepsilon)^{n} J_{t}^{n} f$, where $J_{t}^{n} f$ is the $n$th chaos in the chaos expansion of $S_{t} f$; hence

$$
\begin{equation*}
E\left[\left(S_{t}^{\varepsilon} f-S_{t} f\right)^{2}\right]=\sum_{n \geq 1}\left(1-(1-\varepsilon)^{2 n}\right) E\left[\left(J_{t}^{n} f\right)^{2}\right] . \tag{9.23}
\end{equation*}
$$

Hence it is clear that the $L^{2}(P)$-limit of $S_{t}^{\varepsilon} f$ as $\varepsilon$ goes to 0 is $S_{t} f$.
Let $\left(X_{t}^{\varepsilon}, Y_{t}^{\varepsilon}\right)$ be the Markov process associated with $P_{t}^{(2), \varepsilon}=E\left[S_{t}^{\varepsilon \otimes 2}\right]$ and let $\psi_{t}^{\varepsilon}=d\left(X_{t}^{\varepsilon}, Y_{t}^{\varepsilon}\right) . \psi_{t}^{\varepsilon}$ is a diffusion with generator $L_{\varepsilon}$. It is easy to see that $L_{\varepsilon}=\left(1-(1-\varepsilon)^{2}\right) L_{1}+(1-\varepsilon)^{2} L$ [note that $A_{\varepsilon}^{(2)}=A \otimes I+I \otimes A+(1-\varepsilon)^{2} C=$ $\left.A_{1}^{(2)}+(1-\varepsilon)^{2}\left(A^{(2)}-A_{1}^{(2)}\right)\right]$, and $L_{\varepsilon}=\sigma_{\varepsilon}^{2}(\varphi) \frac{d^{2}}{d \varphi^{2}}+b_{\varepsilon}(\varphi) \frac{d}{d \varphi}$, with

$$
\begin{align*}
\sigma_{\varepsilon}^{2}(\varphi) & =\left(1-(1-\varepsilon)^{2}\right) \sigma_{1}^{2}(\varphi)+(1-\varepsilon)^{2} \sigma^{2}(\varphi),  \tag{9.24}\\
b_{\varepsilon}(\varphi) & =\left(1-(1-\varepsilon)^{2}\right) b_{1}(\varphi)+(1-\varepsilon)^{2} b(\varphi) . \tag{9.25}
\end{align*}
$$

Let us remark that $L_{1}$ is the generator of the diffusion distance between two independent Brownian motions on $S^{d}$. Note that, as $\varphi$ goes to 0 ,
(9.26) $\sigma_{1}^{2}(\varphi) \sim \sigma_{1}^{2}(0)=2(a+b) K G(0) \quad$ and $\quad b_{1}(\varphi) \sim \frac{2(d-1)}{\varphi}(a+b) K G(0)$ and $\sigma_{\varepsilon}^{2}(\varphi)=\left(1-(1-\varepsilon)^{2}\right) \sigma_{1}^{2}(\varphi)\left(1+O\left(\varphi^{\alpha}\right)\right)$ and $b_{\varepsilon}(\varphi)=\left(1-(1-\varepsilon)^{2}\right) b_{1}(\varphi)(1+$ $\left.O\left(\varphi^{\alpha}\right)\right)$. Studying the scale function $s_{\varepsilon}$ of $\psi_{t}^{\varepsilon}$, we get that $s_{\varepsilon}(0+)=s_{1}(0+)=-\infty$ (as two independent Brownian motions cannot meet each other on $S^{d}$ ). We still have $s_{\varepsilon}(\pi-)=\infty$. Hence $\left.\psi_{t}^{\varepsilon} \in\right] 0, \pi[$ for all positive $t$.

Let $m_{\varepsilon}$ be the speed measure of $\psi_{t}^{\varepsilon}$. Let $g_{\varepsilon}(x)=m_{\varepsilon}(d x) / d x$. As $m_{\varepsilon}(] 0, \pi[)$ $<\infty, m_{\varepsilon}$ is an invariant finite measure for the diffusion $\psi_{t}^{\varepsilon}$. As $\lim _{\varepsilon \rightarrow 0} \sigma_{\varepsilon}^{2}=\sigma^{2}$ and $\lim _{\varepsilon \rightarrow 0} b_{\varepsilon}=b$, we get that $\lim _{\varepsilon \rightarrow 0} g_{\varepsilon}(x)=g(x)$. Let us note $\varepsilon^{\prime}=1-(1-\varepsilon)^{2}$ and let

$$
\begin{equation*}
f\left(\varepsilon^{\prime}, \varphi\right)=\frac{\varepsilon^{\prime} b_{1}(\varphi)+\left(1-\varepsilon^{\prime}\right) b(\varphi)}{\varepsilon^{\prime} \sigma_{1}^{2}(\varphi)+\left(1-\varepsilon^{\prime}\right) \sigma^{2}(\varphi)} \tag{9.27}
\end{equation*}
$$

This function increases with $\varepsilon^{\prime}$ if $\frac{b_{1}(\varphi)}{\sigma_{1}^{2}(\varphi)} \geq \frac{b(\varphi)}{\sigma^{2}(\varphi)}$. As $\frac{b_{1}(\varphi)}{\sigma_{1}^{2}(\varphi)}-\frac{b(\varphi)}{\sigma^{2}(\varphi)} \sim(d-1-\mu) \frac{1}{\varphi}$ as $\varphi$ goes to 0 and as $(d-1-\mu)$ is positive, there exists $\varphi_{0}$ such that, for any $\varphi<\varphi_{0}, f\left(\varepsilon^{\prime}, \varphi\right) \geq \frac{b(\varphi)}{\sigma^{2}(\varphi)}=f(0, \varphi)$ and, for $\varepsilon^{\prime}<1 / 2$,

$$
g_{\varepsilon}(x) \leq \frac{2}{\sigma^{2}(x)} \exp \left(-\int_{x}^{\varphi_{0}} \frac{b(\varphi)}{\sigma^{2}(\varphi)} d \varphi\right) C_{\varphi_{0}}
$$

where $C_{\varphi_{0}}=\sup _{\varepsilon \in[0,1]} \exp \left(\int_{\varphi_{0}}^{x_{0}} f\left(\varepsilon^{\prime}, \varphi\right) d \varphi\right)<\infty$. The Lebesgue dominated convergence theorem implies that $g_{\varepsilon}$ converges in $L^{1}([0, \pi])$ toward $g$.

Let $f$ and $g$ be continuous functions. Then $E\left[f\left(X_{t}^{\varepsilon}\right) g\left(Y_{t}^{\varepsilon}\right)\right]=E\left[S_{t}^{\varepsilon} f(x) \times\right.$ $\left.S_{t}^{\varepsilon} g(y)\right]$. Since $S_{t}^{\varepsilon} f$ and $S_{t}^{\varepsilon} g$ converge respectively toward $S_{t} f$ and $S_{t} g$ when $\varepsilon$ goes to 0 in $L^{2}(P)$, we get that $\left(X_{t}^{\varepsilon}, Y_{t}^{\varepsilon}\right)$ converges in distribution toward $\left(X_{t}, Y_{t}\right)$ when $\varepsilon$ goes to 0 . This also implies that $\psi_{t}^{\varepsilon}$ converges in distribution toward $\psi_{t}$ when $\varepsilon$ goes to 0 .

Since $m_{\varepsilon}$ is an invariant measure, for any continuous function $f$ on $[0, \pi]$, we have

$$
\begin{equation*}
\int E\left[f\left(\psi_{t}^{\varepsilon}\right) \mid \psi_{0}^{\varepsilon}=x\right] m_{\varepsilon}(d x)=\int f d m_{\varepsilon} \tag{9.28}
\end{equation*}
$$

Since

$$
\begin{aligned}
& \left|\int E\left[f\left(\psi_{t}^{\varepsilon}\right) \mid \psi_{0}^{\varepsilon}=x\right] m_{\varepsilon}(d x)-\int E\left[f\left(\psi_{t}\right) \mid \psi_{0}=x\right] g(x) d x\right| \\
& \quad \leq\|f\|_{\infty} \int_{0}^{\pi}\left|g_{\varepsilon}(x)-g(x)\right| d x \\
& \quad+\left|\int_{0}^{\pi}\left(E\left[f\left(\psi_{t}^{\varepsilon}\right) \mid \psi_{0}^{\varepsilon}=x\right]-E\left[f\left(\psi_{t}\right) \mid \psi_{0}=x\right]\right) g(x) d x\right|,
\end{aligned}
$$

we get that (because $g_{\varepsilon}$ converges in $L^{1}([0, \pi])$ toward $g$ and $\psi_{t}^{\varepsilon}$ converges in distribution toward $\psi_{t}$ )

$$
\begin{aligned}
\int E\left[f\left(\psi_{t}\right) \mid \psi_{0}=x\right] m(d x) & =\lim _{\varepsilon \rightarrow 0} \int E\left[f\left(\psi_{t}^{\varepsilon}\right) \mid \psi_{0}^{\varepsilon}=x\right] m_{\varepsilon}(d x) \\
& =\lim _{\varepsilon \rightarrow 0} \int f d m_{\varepsilon}=\int f d m .
\end{aligned}
$$

This implies that $g(x) d x$ is an invariant measure for $\psi_{t}$ and $m(d x)=g(x) d x$. Since $m(] 0, x[)<\infty$ for all $x \in] 0, \pi\left[\right.$, the diffusion $\psi_{t}$ is not absorbed in 0 and is reflected in 0 .

In this case, 0 is a closed regular boundary point. This point is instantaneously reflecting since $m(\{0\})=0$. This implies the existence of a positive $t$, a positive $r$ and $p \in] 0,1]$ such that, for any $x \in] 0, \pi\left[, P_{x}\left[\psi_{t} \geq r\right] \geq p\right.$. Then, applying Proposition 6.8, the statistical solution is not a flow of maps. This completes the proof of Theorem 9.4.

For $\alpha>2$, the statistical solution is an isotropic Brownian flow of diffeomorphisms. In [25], the Lyapunov exponents of this flow are computed. The sign of the first Lyapunov exponent $\lambda_{1}(\alpha, d)$ describes the stability of the flow. It is unstable if $\lambda_{1} \geq 0$ and stable if $\lambda_{1}<0$. The computation of $\lambda_{1}(\alpha, d)$ gives

$$
\begin{align*}
\lambda_{1}= & \frac{(d-4) a+d b}{d+2} \zeta(\alpha-1)+\left(\frac{d-1}{d+2}\right)[(d-4) a+d b] \zeta(\alpha)  \tag{9.29}\\
& -d\left(\frac{2(d-1) a+d b}{d+2}\right) \zeta(\alpha+1),
\end{align*}
$$

where $\zeta(\alpha)=\sum_{l \geq 1} \frac{1}{l^{\alpha}}$ is the zeta function. Therefore, we have $\lambda_{1}(\alpha, d)=0$ if and only if

$$
\begin{align*}
\eta & =\eta(\alpha, d) \\
& =\frac{-(d-4) \zeta(\alpha-1)-(d-1)(d-4) \zeta(\alpha)+2 d(d-1) \zeta(\alpha+1)}{4 \zeta(\alpha-1)+4(d-1) \zeta(\alpha)+d(d-2) \zeta(\alpha+1)} . \tag{9.30}
\end{align*}
$$

It is easy to see that, for fixed $\eta, \lim _{\alpha \rightarrow 2+} \lambda_{1}(\alpha, d)=+\infty$ if $d \geq 4$ or if $\eta>\frac{1}{2}-\frac{d-2}{4}=\frac{4-d}{4}$ and that $\lim _{\alpha \rightarrow 2+} \lambda_{1}(\alpha, d)=-\infty$ if $\eta<\frac{4-d}{4}$. Note that $\lim _{\alpha \rightarrow 2-} 1-\frac{d}{\alpha^{2}}=\lim _{\alpha \rightarrow 2-} \frac{1}{2}-\frac{(d-2)}{2 \alpha}=\frac{4-d}{4}$. This shows that coalescence appears when $\lambda_{1}$ goes to $-\infty$ and splitting appears when $\lambda_{1}$ goes to $+\infty$.

The results of this section are given by phase diagrams in Appendix B.

## 10. Isotropic statistical solution on $\mathbb{R}^{d}$.

10.1. Stationary and isotropic covariance functions on $\mathbb{R}^{d}$. On $\mathbb{R}^{d}$ with $d \geq 2$, the stationary isotropic covariance functions $C$ are (see [20]) such that $C^{i j}(x, y)=$ $C^{i j}(x-y)$, for $(x, y) \in \mathbb{R}^{d} \times \mathbb{R}^{d}$, with

$$
\begin{equation*}
C^{i j}(z)=\delta^{i j} B_{N}(\|z\|)+\frac{z^{i} z^{j}}{\|z\|^{2}}\left(B_{L}(\|z\|)-B_{N}(\|z\|)\right) \tag{10.1}
\end{equation*}
$$

with

$$
\begin{align*}
B_{L}(r)= & \iint \cos \left(\rho u_{1} r\right) u_{1}^{2} \omega(d u)\left(F_{L}(d \rho)-F_{N}(d \rho)\right) \\
& +\iint \cos \left(\rho u_{1} r\right) \omega(d u) F_{N}(d \rho)  \tag{10.2}\\
B_{N}(r)= & \iint \cos \left(\rho u_{1} r\right) u_{2}^{2} \omega(d u)\left(F_{L}(d \rho)-F_{N}(d \rho)\right) \\
& +\iint \cos \left(\rho u_{1} r\right) \omega(d u) F_{N}(d \rho) \tag{10.3}
\end{align*}
$$

$F_{L}$ and $F_{N}$ being finite positive measures on $\mathbb{R}^{+} . \omega(d u)$ is the normalized Lebesgue measure on $S^{d-1} . F_{L}$ and $F_{N}$ represent respectively the gradient part and the zero-divergence part of the associated Gaussian vector field.

For $\alpha$ and $m$ positive reals, let

$$
\begin{aligned}
F(d \rho) & =\frac{\rho^{d-1}}{\left(\rho^{2}+m^{2}\right)^{(d+\alpha) / 2}} d \rho \\
F_{L}(d \rho) & =a F(d \rho) \quad \text { and } \quad F_{N}(d \rho)=\frac{b}{d-1} F(d \rho)
\end{aligned}
$$

where $a$ and $b$ are nonnegative. In the Fourier representation ( $c$ is a positive constant),

$$
\begin{equation*}
\hat{C}^{i j}(k)=c\left(\|k\|^{2}+m^{2}\right)^{-(d+\alpha) / 2}\left(a \frac{k^{i} k^{j}}{\|k\|^{2}}+\frac{b}{d-1}\left(\delta_{i j}-\frac{k^{i} k^{j}}{\|k\|^{2}}\right)\right) \tag{10.4}
\end{equation*}
$$

Notice that, in the Fourier representation, the Laplace operator on vector fields is given by the multiplication by $-\|k\|^{2}$ and the projection $\pi$ on gradient vector fields (in the $L^{2}$-space) by $\frac{k^{i} k^{j}}{\|k\|^{2}}$ [i.e., if $V$ is a vector field and $\hat{V}^{i}(k)$ its Fourier transform, $\left(\pi^{\hat{V}}\right)^{i}(k)=\sum_{j} \frac{k^{i} k^{j}}{\|k\|^{2}} \hat{V}^{j}(k)$.]

Therefore, given an $L^{2}$ vector field, $U^{j}(y)=\int \sum_{i} C^{i j}(x-y) V^{i}(x) d x$ can be expressed as $c\left(-\Delta+m^{2}\right)^{-(d+\alpha) / 2}\left(a \pi V+\frac{b}{d-1}(I-\pi) V\right)$. Since $\langle U, U\rangle_{H}=$ $\langle U, V\rangle_{2}=\int\langle U(x), V(x)\rangle d x$, the self-reproducing space appears to be the $L^{2}$-Sobolev space of order $s=\frac{d+\alpha}{2}$ (defined the same way as in Section 9.1) equipped with the norm

$$
\|V\|^{2}=\frac{1}{a}\|\pi V\|_{s}^{2}+\frac{d-1}{b}\|(I-\pi) V\|_{s}^{2},
$$

where

$$
\|V\|_{s}^{2}=\frac{1}{c}\left\langle\left(-\Delta+m^{2}\right)^{s} V, V\right\rangle_{2} .
$$

Note that if $a$ or $b$ vanishes, the self-reproducing space is $H^{2,(\alpha+d) / 2}$ restricted to divergence-free vector fields or gradient vector fields.
10.2. Phase transitions for the Sobolev statistical solution. Let $P_{t}$ be the semigroup of a Brownian motion on $\mathbb{R}^{d}$ with variance $(a+b) F\left(\mathbb{R}^{+}\right)$. Let $S_{t}$ be the statistical solution associated with $P_{t}$ and $C$. If $\alpha>2, C$ is $C^{2}$. Hence (8.2) is satisfied and the statistical solution $S_{t}$ is a flow of maps.

Suppose $a+b>0$ and let $\eta=\frac{b}{a+b}$. Then we have the following theorem.
Theorem 10.1. For any $\alpha \in] 0,2[$, the following hold:
(a) For $d=2$ or 3 and $\eta<1-\frac{d}{\alpha^{2}}$, the statistical solution is a coalescent flow of maps.
(b) For $d=2$ or 3 and $1-\frac{d}{\alpha^{2}}<\eta<\frac{1}{2}-\frac{(d-2)}{2 \alpha}$, the statistical solution is diffusive with hitting.
(c) For $d=2$ or 3 and $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$ or for $d \geq 4$, the statistical solution is diffusive without hitting.

REMARK. The results of this theorem are exactly the same as for the sphere.

Proof of Theorem 10.1. Let us study the two-point motion $\left(X_{t}, Y_{t}\right)$ starting from $(x, y)$ (with $x \neq y)$. Then $r_{t}=d\left(X_{t}, Y_{t}\right)$ is a diffusion in $\mathbb{R}^{+}$ (eventually reflected in 0), with generator $L=\sigma^{2}(r) \frac{d^{2}}{d r^{2}}+b(r) \frac{d}{d r}$ (see [20]), with

$$
\begin{align*}
\sigma^{2}(r) & =B-B_{L}(r)  \tag{10.5}\\
b(r) & =(d-1) \frac{B-B_{N}(r)}{r} \tag{10.6}
\end{align*}
$$

where $B=B_{L}(0)=B_{N}(0)=\frac{a+b}{d} F\left(\mathbb{R}^{+}\right)$.
LEMMA 10.2. For $\alpha \in] 0,2[$, as $r$ goes to 0 , the following hold:
(i) $\iint \cos \left(\rho u_{1} r\right) \omega(d u) F(d \rho)=F\left(\mathbb{R}^{+}\right)-\alpha_{1} r^{\alpha}+o\left(r^{\alpha}\right)$;
(ii) $\iint \cos \left(\rho u_{1} r\right) u_{1}^{2} \omega(d u) F(d \rho)=\frac{F\left(\mathbb{R}^{+}\right)}{d}-\alpha_{2} r^{\alpha}+o\left(r^{\alpha}\right)$;
(iii) $\iint \cos \left(\rho u_{1} r\right) u_{2}^{2} \omega(d u) F(d \rho)=\frac{F\left(\mathbb{R}^{+}\right)}{d}-\alpha_{3} r^{\alpha}+o\left(r^{\alpha}\right)$;
with $\alpha_{2}=\frac{\alpha+1}{d+\alpha} \alpha_{1}, \alpha_{3}=\frac{1}{d+\alpha} \alpha_{1}$ and

$$
\alpha_{1}=c_{d}\left(\int_{0}^{\infty}(1-\cos x) \frac{d x}{x^{\alpha+1}}\right)\left(\int_{0}^{\pi / 2}(\cos \theta)^{\alpha}(\sin \theta)^{d-2} d \theta\right)
$$

Proof. For $r>0$, making the change of variable $x=\rho u_{1} r$,

$$
\begin{aligned}
\iint & \left(1-\cos \left(\rho u_{1} r\right)\right) \omega(d u) F(d \rho) \\
& =c_{d} \int_{0}^{1} \int_{0}^{\infty}\left(1-\cos \left(\rho u_{1} r\right)\right)\left(1-u_{1}^{2}\right)^{(d-2) / 2} d u_{1} \frac{\rho^{d-1} d \rho}{\left(\rho^{2}+m^{2}\right)^{(d+\alpha) / 2}} \\
& =r^{\alpha} c_{d} \int_{0}^{1}\left(\int_{0}^{\infty}(1-\cos x) \frac{x^{d-1} d x}{\left(x^{2}+r^{2} u_{1}^{2} m^{2}\right)^{(d+\alpha) / 2}}\right) u_{1}^{\alpha}\left(1-u_{1}^{2}\right)^{(d-2) / 2} d u_{1}
\end{aligned}
$$

As

$$
\lim _{r \rightarrow 0} \int_{0}^{\infty}(1-\cos x) \frac{x^{d-1}}{\left(x^{2}+r^{2} u_{1}^{2} m^{2}\right)^{(d+\alpha) / 2}} d x=\int_{0}^{\infty}(1-\cos x) \frac{d x}{x^{\alpha+1}}<\infty
$$

we get that

$$
\begin{aligned}
& \lim _{r \rightarrow 0} \frac{1}{r^{\alpha}} \iint\left(1-\cos \left(\rho u_{1} r\right)\right) \omega(d u) F(d \rho) \\
& \quad=c_{d}\left(\int_{0}^{\infty}(1-\cos x) \frac{d x}{x^{\alpha+1}}\right) I(d-2, \alpha)=\alpha_{1}
\end{aligned}
$$

with $I(n, t)=\int_{0}^{\pi / 2}(\cos \theta)^{t}(\sin \theta)^{n} d \theta=\frac{1}{2} B\left(\frac{n+1}{2}, \frac{t+1}{2}\right)$ for $t \geq 0$ and $n \in \mathbb{N}$, and $B(x, y)=\frac{\Gamma(x) \Gamma(y)}{\Gamma(x+y)}$. This shows (i). Statements (ii) and (iii) can be obtained the
same way with

$$
\alpha_{2}=c_{d} \int_{0}^{\infty}(1-\cos x) \frac{d x}{x^{\alpha+1}} I(d-2, \alpha+2)
$$

and $\alpha_{1}=\alpha_{2}+(d-1) \alpha_{3}$ [note that $\left.\int u_{1}^{2} \omega(d u)=\frac{1}{d}\right]$. It is easy to see that, for $\alpha>0$ and $d \geq 1$,

$$
I(d-2, \alpha+2)=\frac{\alpha+1}{d+\alpha} I(d-2, \alpha) .
$$

Therefore, $\alpha_{2}=\frac{\alpha+1}{d+\alpha} \alpha_{1}$. With the relation $\alpha_{1}=\alpha_{2}+(d-1) \alpha_{3}$, we get that $\alpha_{3}=\frac{1}{d+\alpha} \alpha_{1}$.

REmark 10.3. As $z$ goes to 0 ,

$$
\begin{aligned}
& C^{i j}(z)=B \delta^{i j}-\frac{\alpha_{1}}{d-1}\left[((d-1) a+(d+\alpha-1) b) \delta^{i j}\right. \\
&\left.\quad-\alpha((d-1) a-b) \frac{z^{i} z^{j}}{\|z\|^{2}}\right]\|z\|^{\alpha}(1+o(1)) .
\end{aligned}
$$

Let us note that the dependence on $m$ only appears in $B$.
From this lemma, it is easy to see that, as $r$ goes to 0 ,

$$
\begin{align*}
\sigma^{2}(r) & =\frac{(a+b) \alpha_{1}}{d+\alpha}(\alpha+1-\alpha \eta) r^{\alpha}(1+o(1))  \tag{10.7}\\
b(r) & =\frac{(a+b) \alpha_{1}}{d+\alpha}(d-1+\alpha \eta) r^{\alpha-1}(1+o(1)) \tag{10.8}
\end{align*}
$$

Note that we get the same behavior of $\sigma$ and $b$ around 0 as in Section 9.2.
As in Section 9.2, let us study $s$, the scale function of the diffusion $r_{t}$.
Since $B_{L}(r)$ and $B_{N}(r)$ converge toward 0 as $r$ goes to $\infty$ (as Fourier transforms of finite measures), we get that, as $r$ goes to $\infty, \log \left(s^{\prime}(r)\right) \sim(1-d) \log (r)$. Therefore $s(+\infty)$ is finite if and only if $d \geq 3$.

We also see that $s(0+)=-\infty$ if $\eta>\frac{1}{2}-\frac{(d-2)}{2 \alpha}$ and $s(0+)$ is finite if $\eta<\frac{1}{2}-$ $\frac{(d-2)}{2 \alpha}$.
Let $m$ be the speed measure of the diffusion. Let us study the boundary point 0 .
As $m(] 0, x[)<-\infty$ for any positive $x$ if $\eta>1-\frac{d}{\alpha^{2}}$, as in Section 9.2, with a similar proof, we can prove that if $\eta \in] 1-\frac{d}{\alpha^{2}}, \frac{1}{2}-\frac{(d-2)}{2 \alpha}\left[\right.$, the diffusion $r_{t}$ is instantaneously reflecting at 0 . The only thing there is to change in the proof is to take the test function $f$ in (9.28) with compact support and to remark that $g_{\varepsilon}$ converges toward $g$ in $L_{\text {loc }}^{1}\left(\mathbb{R}^{+}\right)$.

If $\eta<1-\frac{d}{\alpha^{2}}$ (note that $1-\frac{d}{\alpha^{2}} \leq \frac{1}{2}-\frac{d-2}{\alpha}$ ), 0 is an exit boundary point and the diffusion is absorbed by 0 .

Therefore, we get that the following hold:
(i) if $d \geq 3$ and $\eta \in] 1-\frac{d}{\alpha^{2}}, \frac{1}{2}-\frac{(d-2)}{2 \alpha}\left[, r_{t}\right.$ is instantaneously reflecting at 0 and is transient. In this case, as in Section 9.2, $\left(S_{t}\right)_{t \geq 0}$ is diffusive with hitting;
(ii) if $d=2$ and $\eta \in] 1-\frac{d}{\alpha^{2}}, \frac{1}{2}-\frac{(d-2)}{2 \alpha}\left[, r_{t}\right.$ is instantaneously reflecting at 0 and is recurrent. In this case, as in Section 9.2, $\left(S_{t}\right)_{t \geq 0}$ is diffusive with hitting;
(iii) if $d \geq 3$ and $\eta<1-\frac{d}{\alpha^{2}}, r_{t}$ is absorbed at 0 with probability $\frac{s(\infty)-s\left(r_{0}\right)}{s(\infty)-s(0)}$ and converges toward $+\infty$ with probability $\frac{s\left(r_{0}\right)-s(0)}{s(\infty)-s(0)}$. In this case, as in Section 9.2, $\left(S_{t}\right)_{t \geq 0}$ is a coalescent flow of maps;
(iv) if $d=2$ and $\eta<1-\frac{d}{\alpha^{2}}, r_{t}$ is absorbed at 0 a.s. In this case, as in Section 9.2, $\left(S_{t}\right)_{t \geq 0}$ is a coalescent flow of maps.

If $\eta>\frac{1}{2}-\frac{d-2}{2 \alpha}$, then we have that $s(0)=-\infty$. In this case, 0 is an entrance boundary point as $\int_{0+}|s(x)| d m(x)<\infty . r_{t}$ is recurrent if $d=2$ and transient if $d \geq 3$. As in Section 9.2, we prove that $\left(S_{t}\right)_{t \geq 0}$ is diffusive without hitting.

For $\alpha>2$, the statistical solution is a stationary isotropic Brownian flow of diffeomorphisms. In [20], the Lyapunov exponents of this flow are computed. The sign of the first Lyapunov exponent $\lambda_{1}(\alpha, d)$ describes the stability of the flow. It is unstable if $\lambda_{1} \geq 0$ and stable if $\lambda_{1}<0$. The computation of $\lambda_{1}(\alpha, d)$ gives (see [20])

$$
\begin{equation*}
\lambda_{1}=\frac{1}{2(d+2)}((d-4) a+d b) \int \rho^{2} F(d \rho) \tag{10.9}
\end{equation*}
$$

Therefore, we have $\lambda_{1}(\alpha, d)=0$ if and only if $d \leq 4$ and

$$
\begin{equation*}
\eta=\eta(d)=\frac{4-d}{4} \tag{10.10}
\end{equation*}
$$

As in Section 9.2, we see that, for fixed $\eta, \lim _{\alpha \rightarrow 2+} \lambda_{1}(\alpha, d)=+\infty$ if $d \geq 4$ or if $\eta>\frac{1}{2}-\frac{d-2}{4}=\frac{4-d}{4}$ and that $\lim _{\alpha \rightarrow 2+} \lambda_{1}(\alpha, d)=-\infty$ if $\eta<\frac{4-d}{4}$. This shows that coalescence appears when $\lambda_{1}$ goes to $-\infty$ and splitting appears when $\lambda_{1}$ goes to $+\infty$.

Note that $\lim _{\alpha \rightarrow 2-} 1-\frac{d}{\alpha^{2}}=\lim _{\alpha \rightarrow 2-\frac{1}{2}-\frac{(d-2)}{2 \alpha}=\frac{4-d}{4} \text {. } . ~ . ~ . ~}^{2}$.
The results of this section are given by phase diagrams in Appendix B.
11. Reflecting flows. Let $D$ be an open convex domain in $\mathbb{R}^{d}$ with $C^{1}$ boundary $\partial D$. Let $d$ be the Euclidean metric in $\mathbb{R}^{d}$. For any $x \in \partial D$, we denote by $n(x)$ the directed inward unit normal vector to $\partial D$.

Let $P_{t}$ be the semigroup of the Brownian motion in $D$ reflected on $\partial D$. $P_{t}$ is associated with the Dirichlet form $(\mathcal{E}, \mathcal{F})$, where $\mathcal{F}=H^{1}(D)=\{f \in$ $\left.L^{2}(D, d x),|\nabla f| \in L^{2}(D, d x)\right\}$ equipped with the form $\frac{1}{2} \int_{D}|\nabla f|^{2} d x$ (see [11], 1.3.2). Let $C(x, y)$ be a covariance function in $D \times D$ such that $C^{i j}(x, x)=\delta^{i j}$ and satisfying (8.1).

We can construct a statistical solution associated with $P_{t}$ and $C$. Let $P_{t}^{(2)}$ be the semigroup of the two-point motion $\left(X_{t}, Y_{t}\right)$. Let $P_{(\cdot,)}^{(2)}$ be the law of the two-point motion.

We know that $X_{t}$ and $Y_{t}$ are two diffusions in $D$ reflected on $\partial D$. Let $\varphi_{t}$ and $\psi_{t}$ denote the local times of $X_{t}$ and $Y_{t}$ on $\partial D$.

Lemma 11.1. For $h(x, y)=d^{2}(x, y), P_{t}^{(2)} h(x, y) \leq h(x, y) e^{C t}$.
Proof. Let us note that

$$
L^{(2)}=A_{x}+A_{y}+\sum_{i, j} C^{i j}(x, y) \partial_{x_{i}} \partial_{x_{j}} .
$$

From (8.1) and the Lipschitz conditions, we get that

$$
L^{(2)} h(x, y) \leq C h(x, y)
$$

Using Tanaka's formula, there exists a martingale $M_{t}$ such that

$$
\begin{align*}
h\left(X_{t}, Y_{t}\right)-h(x, y)= & M_{t}+\int_{0}^{t} L^{(2)} h\left(X_{s}, Y_{s}\right) d s  \tag{11.1}\\
& +\int_{0}^{t}\left\langle\nabla_{x} h\left(X_{s}, Y_{s}\right), n\left(X_{s}\right)\right\rangle d \varphi_{s} \\
& +\int_{0}^{t}\left\langle\nabla_{y} h\left(X_{s}, Y_{s}\right), n\left(Y_{s}\right)\right\rangle d \psi_{s} . \tag{11.2}
\end{align*}
$$

As $\nabla_{x} h(x, y)=2(x-y)$, using the fact that $D$ is convex, we get that, for $x \in \partial D$,

$$
\left\langle\nabla_{x} h(x, y), n(x)\right\rangle<0 .
$$

This implies that

$$
h\left(X_{t}, Y_{t}\right)-h(x, y) \leq M_{t}+C \int_{0}^{t} h\left(X_{s}, Y_{s}\right) d s
$$

Taking the expectation, we get that $P_{t}^{(2)} h(x, y)-h(x, y) \leq C \int_{0}^{t} P_{s}^{(2)} h(x, y) d s$. Hence we have the lemma.

THEOREM 11.2. The statistical solution is a flow of maps.
Proof. This is the same proof as the proof of Theorem 8.2.

## APPENDIX A

Proof of Lemma 9.5. Take $\varphi \in] 0$, $\pi$ [. At first, we are going to prove that $I(\varphi)=\sum_{l \geq 2} \frac{1}{(l-1)^{\alpha+1}}\left|\frac{d}{d \varphi} \gamma_{l}(\cos \varphi)\right|$ is finite. As $\frac{1}{l^{\alpha}}=\int_{0}^{\infty} e^{-l s} s^{\alpha-1} \frac{d s}{\Gamma(\alpha)}$,

$$
\begin{align*}
I(\varphi) & \leq \int_{0}^{\pi} \int_{0}^{\infty} \sum_{l \geq 1}\left[e^{-s}|z(\varphi, \theta)|\right]^{l} \frac{\left|\frac{d}{d \varphi} z(\varphi, \theta)\right|}{|z(\varphi, \theta)|} s^{\alpha-1} \sin ^{d} \theta \frac{d s d \theta}{\Gamma(\alpha) c_{d}}  \tag{A.1}\\
& \leq \int_{0}^{\pi} \int_{0}^{\infty} f_{\varphi, \theta}(s) d s d \theta=2 \int_{0}^{\pi / 2} \int_{0}^{\infty} f_{\varphi, \theta}(s) d s d \theta
\end{align*}
$$

with

$$
f_{\varphi, \theta}(s)=\frac{e^{-s}\left|\frac{d}{d \varphi} z(\varphi, \theta)\right|}{1-e^{-s}|z(\varphi, \theta)|} \frac{s^{\alpha-1}}{\Gamma(\alpha) c_{d}}
$$

It is easy to see that

$$
\begin{equation*}
\int_{1}^{\infty} f_{\varphi, \theta}(s) d s \leq \frac{1}{\Gamma(\alpha) c_{d}} \int_{1}^{\infty} \frac{e^{-s} s^{\alpha-1}}{\left(1-e^{-s}\right)} d s<\infty \tag{A.3}
\end{equation*}
$$

On the other hand,

$$
\begin{equation*}
\int_{0}^{1} f_{\varphi, \theta}(s) d s \leq \frac{1}{\Gamma(\alpha) c_{d}} \int_{0}^{1} \frac{d s}{1-e^{-s}|z(\varphi, \theta)|}=\frac{1}{\Gamma(\alpha) c_{d}} F_{\varphi}(\theta) \tag{A.4}
\end{equation*}
$$

Let $x_{\varphi}(\theta)=-\log |z(\varphi, \theta)|$, then $F_{\varphi}(\theta)=\int_{x_{\varphi}(\theta)}^{x_{\varphi}(\theta)+1} \frac{d t}{1-e^{-t}} . \operatorname{As} \lim _{\theta \rightarrow 0+} x_{\varphi}(\theta)=0$, we have $F_{\varphi}(\theta) \sim-\log x_{\varphi}(\theta)$ as $\theta$ goes to 0 . From this, we see that $F_{\varphi}(\theta)=$ $O(\log \theta)$ as $\theta$ goes to 0 . This implies that $I(\varphi)$ is finite.

Now, applying the derivation under the integral theorem, we prove that $G$ is differentiable on $] 0, \pi[$ and that, for $\varphi \in] 0, \pi[$,

$$
\begin{equation*}
G^{\prime}(\varphi)=\sum_{l \geq 1} \int_{0}^{\pi} \frac{[z(\varphi, \theta)]^{l-1} \frac{d}{d \varphi} z(\varphi, \theta)}{l^{\alpha}} \sin ^{d} \theta \frac{d \theta}{c_{d}} \tag{A.5}
\end{equation*}
$$

$$
\begin{align*}
& =\int_{0}^{\pi} \int_{0}^{\infty} \sum_{l \geq 1}\left[e^{-s} z(\varphi, \theta)\right]^{l} \frac{\frac{d}{d \varphi} z(\varphi, \theta)}{z(\varphi, \theta)} s^{\alpha-1} \sin ^{d} \theta \frac{d s d \theta}{\Gamma(\alpha) c_{d}}  \tag{A.6}\\
& =\int_{0}^{\pi} \int_{0}^{\infty} \frac{e^{-s} \frac{d}{d \varphi} z(\varphi, \theta)}{1-e^{-s} z(\varphi, \theta)} s^{\alpha-1} \sin ^{d} \theta \frac{d s d \theta}{\Gamma(\alpha) c_{d}} \tag{A.7}
\end{align*}
$$

As $z(\varphi, \pi-\theta)=\overline{z(\varphi, \theta)}$,

$$
G^{\prime}(\varphi)=-\int_{0}^{\pi / 2} \int_{0}^{\infty} \frac{a(s, \varphi)-\sin ^{2} \theta}{b(s, \varphi)+\cos ^{2} \theta} \frac{\cos \varphi}{\sin \varphi} s^{\alpha-1} \sin ^{d} \theta \frac{2 d s d \theta}{\Gamma(\alpha) c_{d}}
$$

with $a(s, \varphi)=\frac{1}{e^{-s} \cos \varphi}$ and $b(s, \varphi)=\frac{\left(1-e^{-s} \cos \varphi\right)^{2}}{e^{-2 s} \sin ^{2} \varphi}$. Changing variables $(s=t \varphi)$,

$$
\begin{align*}
-\frac{G^{\prime}(\varphi)}{\varphi^{\alpha-1}} & =\int_{0}^{\pi / 2} \int_{0}^{\infty} \frac{a(t \varphi, \varphi)-\sin ^{2} \theta}{b(t \varphi, \varphi)+\cos ^{2} \theta} \frac{\varphi \cos \varphi}{\sin \varphi} t^{\alpha-1} \sin ^{d} \theta \frac{2 d t d \theta}{\Gamma(\alpha) c_{d}}  \tag{A.8}\\
& =\int_{0}^{\pi / 2} \int_{0}^{\infty} I(t, \varphi, \theta) d t d \theta \tag{A.9}
\end{align*}
$$

Let $\varepsilon>0$. There exists a positive constant $C_{\varepsilon}$ such that, for any $t \in[0, \varepsilon]$,

$$
\begin{equation*}
0 \leq I(t, \varphi, \theta) \leq C_{\varepsilon} t^{\alpha-1} \tag{A.10}
\end{equation*}
$$

Note also that

$$
\begin{equation*}
I(t, \varphi, \theta) \leq C_{d, \alpha} \frac{t^{2} \varphi^{2} e^{-t \varphi}}{\left(1-e^{-t \varphi}\right)^{2}} t^{\alpha-3} \tag{A.11}
\end{equation*}
$$

where $C_{d, \alpha}$ is a positive constant. Let $C=C_{d, \alpha} \sup _{x>0} \frac{x^{2} e^{-x}}{\left(1-e^{-x}\right)^{2}}<\infty$. Then, for any positive $t$,

$$
\begin{equation*}
0 \leq I(t, \varphi, \theta) \leq C t^{\alpha-3} \tag{A.12}
\end{equation*}
$$

As $F(t)=C_{\varepsilon} t^{\alpha-1} \mathbb{1}_{0<t \leq \varepsilon}+C t^{\alpha-3} \mathbb{1}_{t>\varepsilon}$ belongs to $L^{1}(d \theta \otimes d t)$ for $\left.\alpha \in\right] 0,2[$, $\lim _{\varphi \rightarrow 0} a(t \varphi, \varphi)=1$ and $\lim _{\varphi \rightarrow 0} b(t \varphi, \varphi)=t^{2}$, by the Lebesgue dominated convergence theorem,
(A.13) $\lim _{\varphi \rightarrow 0} \frac{G^{\prime}(\varphi)}{\varphi^{\alpha-1}}=-\int_{0}^{\pi / 2} \int_{0}^{\infty} \frac{\cos ^{2} \theta}{t^{2}+\cos ^{2} \theta} t^{\alpha-1} \sin ^{d} \theta \frac{2 d \theta d t}{c_{d} \Gamma(\alpha)}=-\alpha K$.

We have proved the second limit. The first limit is easy to obtain as

$$
\begin{aligned}
G(0)-G(\varphi) & =-\int_{0}^{\varphi} G^{\prime}(x) d x \\
& =-K \varphi^{\alpha}+o\left(\varphi^{\alpha}\right)
\end{aligned}
$$

This finishes the proof of the lemma.

## APPENDIX B

Phase diagrams for the Sobolev statistical solutions. Figures $1-7$ give results of Sections 9 and 10.

Let us remark that, when $\alpha<2$, the diagrams are exactly the same for the sphere and for the plane. For the sphere, we see that, for $\alpha>2$ and $\eta \leq 2-\frac{\zeta(\alpha)}{\zeta(\alpha+1)}$, the flow becomes stable when $d$ goes to $\infty$ : (9.30) implies that $\lim _{d \rightarrow \infty} \eta(\alpha, d)=$ $2-\frac{\zeta(\alpha)}{\zeta(\alpha+1)}$ for $\alpha>2$. We see that, for any $d$ and $\eta \in[0,1[$, the flow becomes stable when $\alpha$ goes to $\infty$ : (9.30) implies that $\lim _{\alpha \rightarrow \infty} \eta(\alpha, d)=1$.


Fig. 1. Phase diagram on $S^{2}$.


Fig. 2. Phase diagram on $\mathbb{R}^{2}$.


Fig. 3. Phase diagram on $S^{3}$.


Fig. 4. Phase diagram on $\mathbb{R}^{3}$.


Fig. 5. Phase diagram on $S^{4}$.


Fig. 6. Phase diagram on $S^{5}$.


Fig. 7. Phase diagram on $S^{50}$.

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