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# CONSTRUCTION OF SET-VALUED DUAL PROCESSES ON MANIFOLDS 

MARC ARNAUDON, KOLÉHÈ COULIBALY-PASQUIER, AND LAURENT MICLO


#### Abstract

The purpose of this paper is to construct a Brownian motion $X:=\left(X_{t}\right)_{t \geqslant 0}$ taking values in a Riemannian manifold $M$, together with a compact valued process $D:=$ $\left(D_{t}\right)_{t \geqslant 0}$ such that, at least for small enough $\mathscr{F}^{D}$-stopping time $\tau \geqslant 0$ and conditioned to $\mathscr{F}_{\tau}^{D}$, the law of $X_{\tau}$ is the normalized Lebesgue measure on $D_{\tau}$. This intertwining result is a generalization of Pitman theorem. We first construct regular intertwined processes related to Stokes' theorem. Then using several limiting procedures we construct synchronous intertwined, free intertwined, mirror intertwined processes. The local times of the Brownian motion on the (morphological) skeleton or the boundary of $D$ plays an important role. Several example with moving intervals, discs, annulus, symmetric convex sets are investigated. KEYWORDS. Brownian motions on Riemannian manifolds, intertwining relations, setvalued dual processes, couplings of primal and dual processes, stochastic mean curvature evolutions, boundary and skeleton local times, generalized Pitman theorem. MSC2010 primary: 60J60, secondary: 60J65, $60 \mathrm{H} 10,58 \mathrm{~J} 65,53 \mathrm{C} 44,60 \mathrm{~J} 55,35 \mathrm{~K} 93$.


## 1. Introduction and main results

Let $M$ be a $d$-dimensional complete Riemannian manifold. We fix a point $o \in M$ for convenience. Denote respectively by $\rho, \mu$ and $\mu$, the Riemannian distance, the Lebesgue measure on $M$ and the corresponding $(d-1)$-Hausdorff measure. The main objective of this paper is to construct intertwined processes and to solve Conjecture 6 in [6] in the case of Brownian motion $\left(X_{t}\right)_{t \geqslant 0}$ and stochastic modified mean curvature flow $\left(D_{t}\right)_{t \geqslant 0}$. This conjecture says that an intertwined construction in the sense of Definition 1.1 is always possible.
Definition 1.1. We say that a Brownian motion $X=\left(X_{t}\right)_{t \geqslant 0}$ in $M$ and a Markov process $D=\left(D_{t}\right)_{t \geqslant 0}$, with values in subsets of $M$ and continuous with respect to the Hausdorff topology, are $\tau$-intertwined where $\tau$ is a positive stopping time in the filtration $\mathscr{F}^{D}$ of $D$ if for all bounded $\mathscr{F}^{D}$-stopping time $\tau^{\prime}$ smaller than $\tau$, conditioned on $\mathscr{F}_{\tau^{\prime}}^{D}, X_{\tau^{\prime}}$ has uniform law in $D_{\tau^{\prime}}$ (and in particular $X_{\tau^{\prime}} \in \bar{D}_{\tau^{\prime}}$ ). We say that $X$ and $D$ are intertwined if they are $\tau$-intertwined, $\tau$ being the lifetime of $(X, D)$, assumed to be a.s. positive and $\mathscr{F}^{D}$-measurable.

This is a generic definition, below stronger topologies on subsets of $M$ will be considered.

Our main results are Theorems 2.8, 3.5 and 4.1 presenting such joint constructions of the primal Brownian motion $\left(X_{t}\right)_{t \geqslant 0}$ and the dual domain-valued $\left(D_{t}\right)_{t \geqslant 0}$ processes. The coupling of Theorem 2.8 consists in the infinite-dimensional stochastic differential equation (2.10), based on a function $f:(x, D) \mapsto f(x, D)$ which is a deformation of the signed distance from $x \in M$ to the boundary of the domain $D$ (see Assumption (2.2) for

[^0]the precise requirements). Theorem 3.5 is obtained by specifying some approximating functions $f$. Given the trajectory $\left(X_{t}\right)_{t \geqslant 0}$ of the Brownian motion, we construct the domain evolution $\left(D_{t}\right)_{t \geqslant 0}$ using the local time of $\left(X_{t}\right)_{t \geqslant 0}$ on the skeletons of $\left(D_{t}\right)_{t \geqslant 0}$ and the mean curvatures of the normal foliations of these domains (see (3.30)). Other approximating functions $f$ lead to Theorem 4.1, where the prominent role is played by the local time at the boundary. This situation is in some sense opposite to the previous one, since the driving Brownian motion is now independent from $\left(X_{t}\right)_{t \geqslant 0}$, while it was as correlated as it can be in Theorem 3.5. These theoretical results are illustrated by the fundamental examples of Section 5. First we recover the intertwining relation between the real Brownian motion and the three-dimensional Bessel process. Next we deal with rotationally symmetric manifolds. Finally we present the application of our results to symmetric convex domains in the plane, even if the detailed proofs are deferred to a forthcoming paper.

A first motivation for such constructions comes from the quantitative investigation of convergence to equilibrium of diffusions on manifold. Assume that $X$ and $D$ are $\tau$ intertwined, where $\tau$ is the hitting time by $D$ of the whole state space $M$. If furthermore $\tau$ is finite (typically true when $M$ is compact), then the Riemannian measure can be normalized into a probability and $\tau$ is a strong stationary time for the Brownian motion $X$, i.e. a stopping time such that $\tau$ and $X_{\tau}$ are independent and $X_{\tau}$ is uniformly distributed. In this situation, the tail distributions of $\tau$ provide quantitative estimates for the speed of convergence of the Brownian motion toward equilibrium, in the separation sense, see Diaconis and Fill [7] for the general description of this approach, in the case of finite state space. This probabilistic method is an alternative to the functional inequality approach, see e.g. the book [3] of Bakry, Gentil and Ledoux, and is based on other geometric considerations, as we will see in the sequel. Other motivations for the couplings of primal and dual processes in the context of diffusions can be found in Machida [11] and [14].

## 2. Intertwined dual processes: EXistence in connection with Stoke's FORMULA

In this section we make a construction of intertwined processes $X$ and $D$ based on the Stokes' Formula (2.1) below. Consider a relatively compact domain $D$ in $M$ with $C^{2}$ boundary. Let $f: \bar{D} \rightarrow \mathbb{R}$ a $C^{2}$ function such that $\left.\nabla f\right|_{\partial D}=N^{D}$ the normal inward vector on boundary. Then by Stoke's formula, for any $C^{2}$ function $g: \bar{D} \rightarrow \mathbb{R}$,

$$
\begin{equation*}
-\int_{\partial D} g d \underline{\mu}=\int_{\partial D} g\left\langle\nabla f,-N^{D}\right\rangle d \underline{\mu}=\int_{D} g \Delta f d \mu+\int_{D}\langle\nabla g, \nabla f\rangle d \mu \tag{2.1}
\end{equation*}
$$

For $\alpha \in(0,1)$, denote by $\mathscr{D}^{2+\alpha}$ the set of relatively compact connected open subsets $D$ of $M$ with $C^{2+\alpha}$ boundary.
Definition 2.1. For a given $\alpha \in(0,1), \varepsilon>0$, we denote by $\mathcal{F}^{\alpha, \varepsilon}$ the set of $D \in \mathscr{D}^{2+\alpha}$ such that

- $D \subset B(o, R)$ the Riemannian ball centered at $o$ with radius $R=1 / \varepsilon$;
- $\rho(\partial D, S(D)) \geqslant \varepsilon$, where $S=S(D)$ is the skeleton of $D$ (see appendix A for details);
- $\rho\left(\partial D, S^{\text {out }}(D)\right) \geqslant \varepsilon$, where $S^{\text {out }}(D)$ is the outer skeleton of $D$, i.e. the skeleton of $(\bar{D})^{c}$.
On $\mathcal{F}^{\alpha, \varepsilon},\left(\widetilde{D}_{t}\right)_{t \in\left[0, \tau_{\varepsilon}\right]}$ will be a diffusion associated to the generator $\widetilde{\mathscr{L}}$ defined in (2.12) and $\tau_{\varepsilon} \in(0,+\infty]$ will be the exiting time from $\mathcal{F}^{\alpha, \varepsilon}$. We extend the trajectory $\left(\widetilde{D}_{t}\right)_{t \in\left[0, \tau_{\varepsilon}\right]}$ by taking $\widetilde{D}_{t}=\widetilde{D}_{\tau_{\varepsilon}}$ for any $t>\tau_{\varepsilon}$. It amounts to imposing that $\widetilde{\mathscr{L}}$ vanishes outside
$\mathcal{F}^{\alpha, \varepsilon}$. It is possible to define in the same way $\left(\widetilde{D}_{t}\right)_{t \in[0, \tau)}$ on $\mathcal{D}^{2+\alpha}$ (which coincides with $\cup_{\varepsilon>0} \mathcal{F}^{\alpha, \varepsilon}$ ), where $\tau$ is the exiting time from $\mathcal{D}^{2+\alpha}$. But it will be more convenient for us to work with a process with an infinite lifetime (to be able to apply Proposition D. 3 in Appendix D) and whose set of values has a boundary which is well-separated from the skeleton.

Let $\beta \in\{0, \alpha\}$. For $D_{0} \in \mathcal{D}^{2+\beta}$ and $\delta>0$ small enough, a $\delta$-neighborhood of $D_{0}$ is defined as follow:

$$
\mathcal{V}_{\delta}^{2+\beta}\left(D_{0}\right):=\left\{\operatorname{int}\left(\exp _{\partial D_{0}}(f)\right), f \in C^{2+\beta}\left(\partial D_{0}\right),\|f\|_{C^{2+\beta}\left(\partial D_{0}\right)}<\delta\right\}
$$

where for $f \in C^{2+\beta}\left(\partial D_{0}\right)$

$$
\exp _{\partial D_{0}}(f):=\left\{\exp _{x}\left(f(x) N^{D_{0}}(x)\right), x \in \partial D_{0}\right\}
$$

( $\exp$ being the exponential map in $M$ ), and $\operatorname{int}\left(\exp _{\partial D_{0}}(f)\right)$ is the interior of the hypersurface $\exp _{\partial D_{0}}(f)$, oriented by the orientation of $D_{0}$. Let $\eta\left(\partial D_{0}\right)>0$ be the radius of the maximal tubular neighborhood of $\partial D_{0}$. Notice that $\delta<\eta\left(\partial D_{0}\right)$ garantees that all elements of $\mathcal{V}_{\delta}^{2+\beta}\left(D_{0}\right)$ are regular deformations of $D_{0}$. Also notice that all elements $D$ of $\mathcal{F}^{\alpha, \varepsilon}$ have $\eta(\partial D) \geqslant \varepsilon$.

We identify two domains $D_{1}, D_{2} \in \mathcal{V}_{\delta}^{2+\beta}\left(D_{0}\right)$ with the functions $f_{1}, f_{2} \in C^{2+\beta}\left(\partial D_{0}\right)$ such that $D_{1}=\operatorname{int}\left\{\exp _{\partial D_{0}}\left(f_{1}\right)\right\}$ and $D_{2}=\operatorname{int}\left\{\exp _{\partial D_{0}}\left(f_{2}\right)\right\}$ and we define a local distance

$$
\begin{equation*}
d_{\beta, D_{0}}\left(D_{1}, D_{2}\right):=\left\|f_{1}-f_{2}\right\|_{C^{2+\beta}\left(\partial D_{0}\right)} . \tag{2.2}
\end{equation*}
$$

## Assumption 2.2.

- The function

$$
\begin{aligned}
f: M \times \mathcal{F}^{\alpha, \varepsilon} & \rightarrow \mathbb{R} \\
(x, D) & \mapsto f(x, D)=f^{D}(x)
\end{aligned}
$$

is a $C^{2+\alpha}$ function in the two variables (the differential in $D$ is in the sense of Fréchet with respect to the above local Banach structure) satisfying

$$
\begin{equation*}
\left\|\nabla f^{D}\right\|_{\infty} \leqslant 1 \tag{2.3}
\end{equation*}
$$

and the functions $f^{D}$ coincide to the signed distance to the boundary $\rho_{\partial D}^{+}$(positive inside $D$ ) in a neighbourhood of $\partial D$. The functions $f^{D}$ have bounded Hessian, uniformly in $D \in \mathcal{F}^{\alpha, \varepsilon}$. Furthermore, we assume that the coefficients of the $\alpha$ Hölderianity of Hess $f^{D}$ are uniformly bounded over $\mathcal{F}^{\alpha, \varepsilon}$.

- There exists a positive integer $m$ and a $C^{1}$ map

$$
\begin{aligned}
\sigma_{c}: M \times \mathcal{F}^{\alpha, \varepsilon} & \rightarrow \Gamma\left(T M \otimes\left(\mathbb{R}^{m}\right)^{*}\right) \\
(x, D) & \mapsto \sigma_{c}(x, D)=\sigma_{c}^{D}(x) \in L\left(\mathbb{R}^{m}, T_{x} M\right)
\end{aligned}
$$

where $\Gamma\left(T M \otimes\left(\mathbb{R}^{m}\right)^{*}\right)$ is the set of sections over $M$ of $T M \otimes\left(\mathbb{R}^{m}\right)^{*}$ and $L\left(\mathbb{R}^{m}, T_{x} M\right)$ is the set of linear maps from $\mathbb{R}^{m}$ to $T_{x} M$, such that the linear map

$$
\begin{align*}
\sigma^{D}(x): \mathbb{R} \times \mathbb{R}^{m} & \rightarrow T_{x} M \\
\left(w_{0}, w\right) & \mapsto w_{0} \nabla f^{D}(x)+\sigma_{\mathrm{c}}^{D}(x)(w) \tag{2.4}
\end{align*}
$$

satisfies

$$
\begin{equation*}
\forall x \in \bar{D}, \quad \sigma^{D}\left(\sigma^{D}\right)^{*}(x)=\operatorname{Id}_{T_{x} M} \tag{2.5}
\end{equation*}
$$

Remark 2.3. The first condition of Assumption 2.2 implies that

$$
\begin{align*}
& \left.\nabla f^{D}\right|_{\partial D}=\left.\left(\nabla \rho_{\partial D}^{+}\right)\right|_{\partial D}\left(=N^{D}\right) \quad \text { and } \\
& \left.\Delta f^{D}\right|_{\partial D}=\left.\left(\Delta \rho_{\partial D}^{+}\right)\right|_{\partial D}\left(=-h^{D}\right) \tag{2.6}
\end{align*}
$$

where $h^{D}$ stands for the mean curvature on $\partial D$. It also implies that the functions $f^{D}$ are uniformly Lipschitz and have uniformly bounded Laplacian. Also, for fixed $x \in \partial D$, varying $D$ successively along a field $K$ normal to the boundary $\partial D$ and along $N^{D}$ for the second derivative:

$$
\begin{align*}
\langle\nabla f(x, \cdot), K\rangle(x) & =-\left\langle N^{D}(x), K(x)\right\rangle \quad \text { and } \\
\nabla d f(x, \cdot)\left(N^{D}, N^{D}\right) & =0 \tag{2.7}
\end{align*}
$$

where $\nabla d f(x, \cdot)$ is the Hessian of $f$ in the second variable.
The second condition of Assumption 2.2 implies that for all $u \in T_{x} M$,

$$
\begin{equation*}
\|u\|^{2}=\left\langle u, \nabla f^{D}(x)\right\rangle^{2}+\sum_{i=1}^{m}\left\langle u, \sigma_{\mathrm{c}}^{D}(x)\left(e_{i}\right)\right\rangle^{2} \tag{2.8}
\end{equation*}
$$

for $e_{1}, \ldots, e_{m}$ an orthonormal basis of $\mathbb{R}^{m}$. In particular, if $x \in \partial D$, taking $u=\nabla f^{D}(x)=$ $N^{D}(x)$, we get since $\left\|N^{D}(x)\right\|=1$ :

$$
\begin{equation*}
0=\left\langle\nabla f^{D}(x), \sigma(x)\left(e_{i}\right)\right\rangle, \quad i=1, \ldots m \tag{2.9}
\end{equation*}
$$

Proposition 2.4. Assumption 2.2 can always be realized, with any $\alpha \in(0,1)$ and $\varepsilon>0$.
Proof. We begin with remarking that for $D \in \mathcal{F}^{\alpha, \varepsilon}, \rho(\partial D, S(D)) \geqslant \varepsilon$. In particular, the distance to $\partial D$ is $C^{2+\alpha}$ on $D_{\varepsilon}:=\{x \in D, \rho(x, \partial D)<\varepsilon\}$. Let $h_{\varepsilon}$ be a smooth nondecreasing function from $[0, \infty)$ to $\mathbb{R}_{+}$such that $h_{\varepsilon}(r)=r$ for $r \in[0, \varepsilon / 2], h_{\varepsilon}(r)=$ $(3 / 4) \varepsilon$ for $r \geqslant \varepsilon$ and $\left\|h_{\varepsilon}^{\prime}\right\|_{\infty} \leqslant 1$. Then $f^{D}:=h_{\varepsilon} \circ \rho_{\partial D}^{+}$satisfies all the requirements of the first condition of Assumption 2.2. Then for constructing $\sigma_{c}^{D}$ we proceed as in [2], Proposition 3.2 taking $\sigma_{1}=\nabla f^{D}$. The wanted regularity in $D$ is easily checked.

Let $W_{t}$ and $W_{t}^{m}$ two independent Brownian motions with values respectively in $\mathbb{R}$ and $\mathbb{R}^{m}$.

The equation we are interested in writes in Itô form for all $y \in \partial D_{t}$ :

$$
\left\{\begin{align*}
d X_{t} & =\left(\nabla f^{D_{t}}\left(X_{t}\right) d W_{t}+\sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right)  \tag{2.10}\\
d \partial D_{t}(y) & =N^{D_{t}}(y)\left(d W_{t}+\left(\frac{1}{2} h^{D_{t}}(y)+\Delta f^{D_{t}}\left(X_{t}\right)\right) d t\right)
\end{align*}\right.
$$

started at a relatively compact domain $D_{0}$ with $C^{2+\alpha}$ boundary and $X_{0}$ such that $\mathscr{L}\left(X_{0}\right)=$ $\mathscr{U}\left(D_{0}\right)$, where $\mathscr{U}\left(D_{0}\right)$ is the uniform probability measure on $D_{0}$. In fact, as in Definition 2.1 , the evolution equation (2.10) is implicitly considered only up to the exit time $\tau_{\varepsilon}$ of $\mathcal{F}^{\alpha, \varepsilon}$ for some fixed $\alpha \in(0,1), \varepsilon>0$, after which the process is assumed not to move.

In (2.10), the processes $\left(D_{t}\right)_{t \geqslant 0}$ and $\left(X_{t}\right)_{t \geqslant 0}$ are fully interacting, since the evolution of one of them depends on the other one. In particular, they are not Markovian by themselves in general.

Another subset-valued process $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ will be interesting for our purposes. It is solution to the evolution equation

$$
\begin{equation*}
d \partial \widetilde{D}_{t}(y)=N^{\widetilde{D}_{t}}(y)\left(d \widetilde{W}_{t}+\left(\frac{1}{2} h^{\widetilde{D}_{t}}(y)-\frac{\underline{\mu}^{\partial \widetilde{D}_{t}}\left(\partial \widetilde{D}_{t}\right)}{\mu\left(\widetilde{D}_{t}\right)}\right) d t\right), \quad \forall y \in \partial \widetilde{D}_{t} \tag{2.11}
\end{equation*}
$$

where $\widetilde{W}_{t}$ is a real-valued Brownian motion.
Notice that the equation for $\widetilde{D}_{t}$ does no longer depend of $X_{t}$, so if the solution is unique, $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ will be Markovian. It is Equation (51) in [6]. Theorem 39 of [6] proves local existence of a solution.

Theorem 2.5. Fix $\alpha \in(0,1)$ and $\varepsilon>0$. Then (2.11) admits a unique global solution. In particular the process $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ is Markovian.

Proof. The proof is a consequence of Theorem 9 in [6]. It can be found in Appendix C.
To describe the generator $\widetilde{\mathscr{L}}$ of $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ we must introduce the following notations. For any smooth function $k$ on $M$, consider the mapping $F_{k}$ on $\mathcal{D}^{2+\alpha}$ by

$$
\forall D \in \mathcal{D}^{2+\alpha}, \quad F_{k}(D):=\int_{D} k d \mu
$$

For any $k, g \in \mathcal{C}^{\infty}(M)$ and any $D \in \mathcal{D}^{2+\alpha}$, define

$$
\begin{align*}
\widetilde{L}\left[F_{k}\right](D) & :=\underline{\mu}^{\partial D}(k) \frac{\underline{\mu}^{\partial D}(\partial D)}{\mu(D)}-\frac{1}{2} \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)  \tag{2.12}\\
\Gamma_{\widetilde{\mathscr{L}}}\left[F_{k}, F_{g}\right](D) & :=\int_{\partial D} k d \underline{\mu} \int_{\partial D} g d \underline{\mu} \tag{2.13}
\end{align*}
$$

Next consider $\mathfrak{A}$ the algebra consisting of the functionals of the form $\mathfrak{F}:=\mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right)$, where $n \in \mathbb{Z}_{+}, k_{1}, \ldots, k_{n} \in \mathcal{C}^{\infty}(M)$ and $\mathfrak{f}: \mathcal{R} \rightarrow \mathbb{R}$ is a $\mathcal{C}^{\infty}$ mapping, with $\mathcal{R}$ an open subset of $\mathbb{R}^{n}$ containing the image of $\mathcal{D}^{2+\alpha}$ by $\left(F_{k_{1}}, \ldots, F_{k_{n}}\right)$. For such a functional $\mathfrak{F}$, define

$$
\begin{align*}
\widetilde{\mathscr{L}}[\widetilde{F}]:= & \sum_{l=1}^{n} \partial_{j} \mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right) \widetilde{\mathscr{L}}\left[F_{k_{l}}\right]  \tag{2.14}\\
& +\sum_{j, l \in \llbracket 1, n \rrbracket}^{n} \partial_{j, l} \mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right) \Gamma_{\widetilde{\mathscr{L}}}\left[F_{k_{j}}, F_{k_{l}}\right]
\end{align*}
$$

To two elements of $\mathfrak{A}, \mathfrak{F}:=\mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right)$ and $\mathfrak{G}:=\mathfrak{g}\left(F_{g_{1}}, \ldots, F_{g_{m}}\right)$, we also associate

$$
\begin{equation*}
:=\sum_{l \in \llbracket n \rrbracket, j \in \llbracket m \rrbracket} \partial_{l} \mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right) \partial_{j} \mathfrak{g}\left(F_{g_{1}}, \ldots, F_{g_{m}}\right) \Gamma_{\widetilde{\mathscr{L}}}\left[F_{k_{l}}, F_{g_{j}}\right] \tag{2.15}
\end{equation*}
$$

Remark 2.6. To see that the above definitions are non-ambiguous, since a priori they could depend on the writing of $\mathfrak{F} \in \mathfrak{A}$ under the form $\mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right)$ and similarly for $\mathfrak{G}$, see Remark 2 of [6]. More generally, the forms of (2.14) and (2.15) are consequences of the diffusion feature of $\widetilde{\mathscr{L}}$, for more on the subject, see e.g. the book of Bakry, Gentil and Ledoux [3].
Remark 2.7. In the above considerations, $\widetilde{\mathscr{L}}$ was defined on $\mathcal{D}^{2+\alpha}$, but from now on, $\widetilde{\mathscr{L}}$ will stand for the restriction of this generator to $\mathcal{F}^{\alpha, \varepsilon}$ and will be zero on $\mathcal{D}^{2+\alpha} \backslash \mathcal{F}^{\alpha, \varepsilon}$, in accordance with Definition 2.1. Similarly, all stochastic differential equations will be valid only up to the stopping time $\tau_{\varepsilon}$.

The interest of Assumption 2.2 comes from the following result:
Theorem 2.8. Let $(x, D) \mapsto f^{D}(x)$ and $(x, D) \mapsto \sigma_{\mathrm{c}}^{D}(x)$ satisfy Assumption 2.2. Then equation (2.10) has a solution $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ started at $D_{0} \in \mathscr{F}^{\alpha, \varepsilon}, X_{0} \sim \mathscr{U}\left(D_{0}\right)$. Moreover the processes $\left(X_{t}\right)_{t \geqslant 0}$ and $\left(D_{t}\right)_{t \geqslant 0}$ are $\tau_{\varepsilon}$-intertwined.

Proof. We prove here the existence of solution to equation (2.10). The intertwining will be a consequence of Proposition 2.11 below.

We begin to prove the existence of a diffusion with modified drift, and then we will get the result by change of probability. The modified equation writes

$$
\left\{\begin{align*}
d \partial D_{t}(y)= & N^{D_{t}}(y)\left(d \widehat{W}_{t}+\left(\frac{1}{2} h^{D_{t}}(y)-\frac{\mu^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}\right) d t\right)  \tag{2.16}\\
d X_{t}= & \left(\nabla f^{D_{t}}\left(X_{t}\right)\left[d \widehat{W}_{t}-\left(\frac{\mu^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}+\Delta f^{D_{t}}\left(X_{t}\right)\right) d t\right]\right. \\
& \left.+\sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right)
\end{align*}\right.
$$

for $\widehat{W}_{t}$ and $W_{t}^{m}$ independent Brownian motions. Notice that the first equation is the same as (2.11). Thus due to Assumption 2.5, $\left(D_{t}\right)_{t \geqslant 0}$ is a diffusion process with generator $\widetilde{\mathscr{L}}$. Then given $D_{t}$, the equation for $X_{t}$

$$
\begin{align*}
d X_{t} & =\left(\nabla f^{D_{t}}\left(X_{t}\right)\left[d \widehat{W}_{t}-\left(\frac{\underline{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}+\Delta f^{D_{t}}\left(X_{t}\right)\right) d t\right]\right. \\
& \left.+\sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right) \tag{2.17}
\end{align*}
$$

can also be solved, since the coefficients in front of $d \widehat{W}_{t}$ and $d W_{t}^{m}$ are Lipschitz, $\sigma^{D}\left(\sigma^{D}\right)^{*}(x)=\operatorname{Id}_{T_{x} M}$ and $\Delta f^{D}$ is bounded and uniformly Hölder continuous (due to Assumption 2.2). Notice that $X_{t}$ remains in $D_{t}$, since when $X_{t} \in \partial D_{t}$, we have, using (2.9) which yields on boundary $\left\langle N^{D_{t}}\left(X_{t}\right), \sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right\rangle=0$,

$$
\begin{align*}
& d\left(\rho_{\partial D_{t}}^{+}\left(X_{t}\right)\right) \\
& =\left\langle\nabla \rho_{\partial D_{t}}^{+}, d X_{t}\right\rangle-\frac{1}{2} h^{D_{t}}\left(X_{t}\right) d t-\left\langle d \partial D_{t}\left(X_{t}\right), N^{D_{t}}\left(X_{t}\right)\right\rangle  \tag{2.18}\\
& =\left\langle N^{D_{t}}\left(X_{t}\right), d X_{t}\right\rangle-\frac{1}{2} h^{D_{t}}\left(X_{t}\right) d t-\left\langle d \partial D_{t}\left(X_{t}\right), N^{D_{t}}\left(X_{t}\right)\right\rangle=0 .
\end{align*}
$$

where we used (2.16) and (2.6). We also have no covariation since the martingale part of $d \partial D_{t}$ acts on the normal flow only, and any normal flow

$$
r \mapsto D(r):=\left\{x \in M, \rho^{+}(x) \geqslant r\right\}
$$

satisfies $\rho_{\partial D(r)}^{+}(x)=\rho_{\partial D(0)}^{+}(x)-r$ for $x \in D(0)$ and $|r|$ small, (see Appendix A).
Once we have a solution to (2.16), make by Girsanov theorem a change of probability such that $\left(W_{t}, W_{t}^{m}\right)$ is a Brownian motion where

$$
\begin{equation*}
W_{t}:=\widehat{W}_{t}-\int_{0}^{t}\left(\frac{\underline{\mu}^{\partial D_{s}}\left(\partial D_{s}\right)}{\mu\left(D_{s}\right)}+\Delta f^{D_{s}}\left(X_{s}\right)\right) d s \tag{2.19}
\end{equation*}
$$

We get a solution to (2.10) in the new probability.
Proposition 2.9. Let $D_{t}$ satisfy

$$
\begin{equation*}
d \partial D_{t}(y)=N^{D_{t}}(y)\left(d W_{t}+\left(\frac{1}{2} h^{D_{t}}(y)+b_{t}\right) d t\right), \quad \forall y \in \partial D_{t} \tag{2.20}
\end{equation*}
$$

for some Brownian motion $W_{t}$ and some adapted locally bounded real-valued process $b_{t}$. Let $\mu_{t}=\mu^{D_{t}}$ be the Lebesgue measure on $D_{t}$ and $\bar{\mu}_{t}=\bar{\mu}^{D_{t}}=\mathscr{U}\left(D_{t}\right)=\frac{\mu^{D_{t}}}{\mu\left(D_{t}\right)}$. Denote
by $\underline{\mu}_{t}=\underline{\mu}^{\partial D_{t}}$ the Lebesgue measure on $\partial D_{t}$ and $\bar{\mu}_{t}=\overline{\bar{\mu}}^{\partial D_{t}}=\frac{\underline{\mu}^{\partial D_{t}}}{\mu\left(D_{t}\right)}$. Let $k$ be a smooth function of $M$. Then

$$
\begin{equation*}
d \mu_{t}(k)=-\underline{\mu}_{t}(k) d W_{t}-\frac{1}{2}\left(2 b_{t} \underline{\mu}_{t}(k)+\underline{\mu}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right)\right) d t \tag{2.21}
\end{equation*}
$$

and

$$
\begin{align*}
d \bar{\mu}_{t}(k) & =\left(-\overline{\underline{\mu}}_{t}(k)+\bar{\mu}_{t}(k) \overline{\underline{\mu}}_{t}\left(\partial D_{t}\right)\right) d W_{t}-\frac{1}{2} \overline{\underline{\mu}}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right) d t  \tag{2.22}\\
& +\left(\overline{\underline{\mu}}_{t}\left(\partial D_{t}\right)+b_{t}\right)\left(-\bar{\mu}_{t}(k)+\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right)\right) d t
\end{align*}
$$

In particular, if $b_{t}=-\underline{\mu}_{t}\left(\partial D_{t}\right)$ we get

$$
\begin{equation*}
d \bar{\mu}_{t}(k)=\left(-\overline{\bar{\mu}}_{t}(k)+\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right)\right) d W_{t}-\frac{1}{2} \overline{\bar{\mu}}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right) d t . \tag{2.23}
\end{equation*}
$$

Proof. For $f$ a smooth function on $M$ and $r \in \mathbb{R}$ sufficiently close to 0 so that $\partial D(r)$ (defined in (A.3) and (A.4)) is a smooth manifold without boundary, let

$$
\begin{equation*}
F(r, k)=\int_{D(r)} k d \mu \tag{2.24}
\end{equation*}
$$

We have

$$
\begin{equation*}
F(r, k)=\int_{\partial D}\left(\int_{r}^{\tau(y)} k(\psi(s)(y)) e^{-\int_{0}^{s} h^{D_{t}}(\psi(u)(y)) d u} d s\right) \underline{\mu}(d y) \tag{2.25}
\end{equation*}
$$

with $\tau(y)$ the hitting time of $S(D)$ by the inward normal flow started at $y$ (defined in (A.1)) and $\psi(s)(y)=\psi(0, s)(y)=\exp _{y}\left(s N_{y}\right)$ defined in (A.5). With this formulation we can differentiate with respect to $r$, to obtain

$$
\begin{equation*}
F^{\prime}(r, k)=-\int_{\partial D} k(\psi(r, y)) e^{-\int_{0}^{r} h^{D_{t}}(\psi(s)(y) d s} \mu(d y) \tag{2.26}
\end{equation*}
$$

Differentiating again we get

$$
\begin{equation*}
F^{\prime \prime}(r, k)=-\int_{\partial D}\left(\left\langle d k, \partial_{r} \psi(r, y)\right\rangle-(k h)(\psi(r, y))\right) e^{-\int_{0}^{r} h^{D_{t}}(\psi(s)(y)) d s} \underline{\mu}(d y) \tag{2.27}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
F^{\prime}(0, k)=-\underline{\mu}(k) \quad \text { and } \quad F^{\prime \prime}(0, k)=\underline{\mu}(k h-\langle d k, N\rangle) . \tag{2.28}
\end{equation*}
$$

This allows us to compute

$$
\begin{equation*}
d\left(F\left(W_{t}, k\right)\right)=F^{\prime}\left(W_{t}, k\right) d W_{t}+\frac{1}{2} F^{\prime \prime}\left(W_{t}, k\right) d t \tag{2.29}
\end{equation*}
$$

and then, since $d W_{t}$ and $\left\langle d \partial D_{t}, N^{D_{t}}\right\rangle(\cdot)$ differ only by a finite variation process

$$
\begin{equation*}
d \mu_{t}(k)=\int_{\partial D_{t}}-k(y)\left\langle d \partial D_{t}(y), N^{D_{t}}(y)\right\rangle+\frac{1}{2}\left(k h^{D_{t}}-\left\langle d k, N^{D_{t}}\right\rangle\right)(y) \underline{\mu}_{t}(d y) \tag{2.30}
\end{equation*}
$$

This yields

$$
\begin{equation*}
d \mu_{t}(k)=\int_{\partial D_{t}} k(y)\left(-d W_{t}-b_{t} d t\right)-\frac{1}{2}\left\langle d k, N^{D_{t}}\right\rangle(y) \underline{\mu}_{t}(d y) d t \tag{2.31}
\end{equation*}
$$

which gives (2.21). In particular, taking $f \equiv 1$ we obtain

$$
\begin{equation*}
d \mu\left(D_{t}\right)=\underline{\mu}_{t}\left(\partial D_{t}\right)\left(-d W_{t}-b_{t} d t\right) \tag{2.32}
\end{equation*}
$$

Now we can compute

$$
\begin{aligned}
& d \bar{\mu}_{t}(k) \\
& =d\left(\frac{\mu_{t}(k)}{\mu\left(D_{t}\right)}\right) \\
& =\frac{1}{\mu\left(D_{t}\right)} d \mu_{t}(k)-\frac{\mu_{t}(k)}{\mu\left(D_{t}\right)^{2}} d \mu\left(D_{t}\right)+\frac{\mu_{t}(k)}{\mu\left(D_{t}\right)^{3}} d\left\langle\mu\left(D_{.}\right)\right\rangle_{t}-\frac{1}{\mu\left(D_{t}\right)^{2}} d\langle\mu .(k), \mu(D .)\rangle_{t} \\
& =\frac{1}{\mu\left(D_{t}\right)} d \mu_{t}(k)-\frac{\mu_{t}(k)}{\mu\left(D_{t}\right)^{2}} d \mu\left(D_{t}\right)+\frac{\mu_{t}(k)}{\mu\left(D_{t}\right)^{3}} \underline{\mu}\left(\partial D_{t}\right)^{2} d t-\frac{1}{\mu\left(D_{t}\right)^{2}} \underline{\mu}_{t}(k) \underline{\mu}_{t}\left(\partial D_{t}\right) d t \\
& =-\overline{\bar{\mu}}_{t}(k)\left(d W_{t}+b_{t} d t\right)-\frac{1}{2} \bar{\mu}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right) d t+\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right)\left(d W_{t}+b_{t} d t\right) \\
& +\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right)^{2} d t-\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right) d t
\end{aligned}
$$

This yields (2.22).
Denote $\tau_{\varepsilon}$ the exiting time of $\left(D_{t}\right)_{t \geqslant 0}$ from $\mathcal{F}^{\alpha, \varepsilon}$. As in Definition 2.1, we stop $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ at $\tau_{\varepsilon}$.

Proposition 2.10. Any solution of equation (2.10) stopped at $\tau_{\varepsilon}$ is a Markov process solution to a martingale problem associated to a generator $\mathscr{L}$ acting in the following way: for any $g, k$ smooth functions on $M$ and

$$
\begin{equation*}
F_{k}(D):=\int_{D} k d \mu \tag{2.33}
\end{equation*}
$$

we have for $(x, D) \in M \times \mathcal{F}^{\alpha, \varepsilon}$,

$$
\begin{align*}
\mathscr{L}\left(g F_{k}\right)(x, D) & =-g(x) \Delta f^{D}(x) \underline{\mu}^{\partial D}(k)-\frac{1}{2} g(x) \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)+\frac{1}{2} F_{k}(D) \Delta g(x)  \tag{2.34}\\
& -\mu^{\partial D}(k)\left\langle\nabla g, \nabla f^{D}\right\rangle(x)
\end{align*}
$$

Proof. From (2.10) and (2.21) with $b_{t}=\Delta f^{D_{t}}\left(X_{t}\right)$ we have

$$
\begin{equation*}
d F_{k}\left(D_{t}\right)=-\underline{\mu}^{\partial D_{t}}(k)\left(d W_{t}+\Delta f^{D_{t}}\left(X_{t}\right) d t\right)-\frac{1}{2} \underline{\mu}^{\partial D_{t}}\left(\left\langle\nabla k, N^{D_{t}}\right\rangle\right) d t \tag{2.35}
\end{equation*}
$$

This implies that

$$
\begin{equation*}
\mathscr{L}\left(F_{k}\right)(x, D)=-\underline{\mu}^{\partial D}(k) \Delta f^{D}(x)-\frac{1}{2} \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right), \tag{2.36}
\end{equation*}
$$

and the covariation of $g\left(X_{t}\right)$ and $F_{k}\left(D_{t}\right)$ is $\Gamma_{\mathscr{L}}\left[g, F_{k}\right]\left(X_{t}, D_{t}\right) d t$ with

$$
\begin{equation*}
\Gamma_{\mathscr{L}}\left[g, F_{k}\right](x, D)=-\underline{\mu}^{\partial D}(k)\left\langle\nabla g, \nabla f^{D}\right\rangle(x) . \tag{2.37}
\end{equation*}
$$

Consequently, using

$$
\begin{equation*}
\mathscr{L}\left(g F_{k}\right)(x, D)=g(x) \mathscr{L}\left(F_{k}\right)(x, D)+F_{k}(D) \frac{1}{2} \Delta g(x)+\Gamma_{\mathscr{L}}\left[g, F_{k}\right](x, D) \tag{2.38}
\end{equation*}
$$

we get (2.34).
It is possible to extend the description of $\mathscr{L}$ to more general functions on $M \times \mathcal{F}^{\alpha, \varepsilon}$ (it vanishes on its complementary set), by replacing $F_{k}$ in (2.34) by a mapping $\mathfrak{F}$ from $\mathfrak{A}$, as presented before Theorem 2.8.

Let $\left(\mathscr{P}_{t}\right)_{t \geqslant 0}$ be the Markovian semi-group associated to the processes $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ solution to (2.10) stopped at $\tau_{\varepsilon}$. This semi-group is associated to $\mathscr{L}$ in the weak sense of martingale problems, as described in Appendix D.

Let $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ be a diffusion process with generator $\widetilde{\mathscr{L}}$ stopped outside $\mathcal{F}^{\alpha, \varepsilon}$, started at $\widetilde{D}_{0}=D_{0}$ (due to Theorem 2.5, this process can be obtained as a solution to the evolution equation (2.11)), $\widetilde{\nu}_{t}$ its law at time $t$ and

$$
\begin{equation*}
\nu_{t}(d D, d x)=\widetilde{\nu}_{t}(d D) \mathscr{U}(D)(d x) \tag{2.39}
\end{equation*}
$$

Proposition 2.11. We have for all smooth functions $g, k$ on $M$ :

$$
\begin{equation*}
\partial_{t} \nu_{t}\left(g F_{k}\right)=\nu_{t}\left(\mathscr{L}\left(g F_{k}\right)\right) \tag{2.40}
\end{equation*}
$$

As a consequence, if $\left(D_{0}, X_{0}\right)$ has law $\nu_{0}$ then for all $t \geqslant 0$, the solution $\left(D_{t}, X_{t}\right)$ to equation (2.10) has law $\nu_{t}$, implying that $\left(X_{t}\right)_{t \geqslant 0}$ and $\left(D_{t}\right)_{t \geqslant 0}$ are $\tau_{\varepsilon}$-intertwined. Moreover $D_{t}$ is a diffusion with generator $\widetilde{\mathscr{L}}$.
Proof. Integrating (2.34) in $x$ with respect to the uniform law $\bar{\mu}^{D}:=\mathscr{U}(D)$ in $D$ yields (2.41)
$-\bar{\mu}^{D}\left(g \Delta f^{D}\right) \underline{\mu}^{\partial D}(k)-\frac{1}{2} \bar{\mu}^{D}(g) \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)+\frac{1}{2} F_{k}(D) \bar{\mu}^{D}(\Delta g)-\underline{\mu}^{\partial D}(k) \bar{\mu}^{D}\left(\left\langle\nabla g, \nabla f^{D}\right\rangle\right)$.
By Stokes theorem,

$$
\begin{equation*}
\bar{\mu}^{D}\left(g \Delta f^{D}+\left\langle\nabla g, \nabla f^{D}\right\rangle\right)=\underline{\bar{\mu}}^{\partial D}\left(g\left\langle\nabla f^{D},-N^{D}\right\rangle\right)=-\underline{\bar{\mu}}^{\partial D}(g) \tag{2.42}
\end{equation*}
$$

so the expression (2.41) writes

$$
\begin{equation*}
H(D):=\underline{\mu}^{\partial D}(k) \underline{\bar{\mu}}^{\partial D}(g)-\frac{1}{2} \bar{\mu}^{D}(g) \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)+\frac{1}{2} F_{k}(D) \bar{\mu}^{D}(\Delta g) \tag{2.43}
\end{equation*}
$$

On the other hand

$$
\begin{equation*}
\nu_{t}\left(g F_{k}\right)=\widetilde{\nu}_{t}\left[\bar{\mu}^{D_{t}}[g] F_{k}\right] \tag{2.44}
\end{equation*}
$$

which implies that

$$
\begin{equation*}
\partial_{t} \nu_{t}\left(g F_{k}\right)=\partial_{t} \widetilde{\nu}_{t}\left(\left(\bar{\mu}^{D_{t}}(g) F_{k}\right)=\widetilde{\nu}_{t}\left(\widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g) F_{k}\right)\right)\right. \tag{2.45}
\end{equation*}
$$

By (2.23),

$$
\begin{equation*}
\widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g)\right)=-\frac{1}{2} \underline{\bar{\mu}}^{\partial D_{t}}\left(\left\langle\nabla g, N^{D_{t}}\right\rangle\right) \tag{2.46}
\end{equation*}
$$

so, taking into account (2.13),

$$
\begin{aligned}
& \widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g) F_{k}\right) \\
& =\bar{\mu}^{D_{t}}(g) \widetilde{\mathscr{L}}\left(F_{k}\right)+F_{k} \widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g)\right)+\Gamma_{\widetilde{\mathscr{L}}}\left[\bar{\mu}^{D_{t}}(g), F_{k}\right] \\
& =\bar{\mu}^{D_{t}}(g)\left\{\underline{\mu}^{\partial D_{t}}(k) \bar{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)-\frac{1}{2} \underline{\mu}^{\partial D_{t}}\left(\left\langle\nabla k, N^{D_{t}}\right\rangle\right)\right\}-\frac{1}{2} \mu^{D_{t}}(k) \bar{\mu}^{\partial D_{t}}\left(\left\langle\nabla g, N^{D_{t}}\right\rangle\right) \\
& -\left(-\bar{\mu}^{\partial D_{t}}(g)+\bar{\mu}^{D_{t}}(g) \bar{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)\right) \underline{\mu}^{\partial D_{t}}(k) \\
& =-\frac{1}{2} \bar{\mu}^{D_{t}}(g) \underline{\mu}^{\partial D_{t}}\left(\left\langle\nabla k, N^{D_{t}}\right\rangle\right)-\frac{1}{2} \mu^{D_{t}}(k) \overline{\underline{\mu}}^{\partial D_{t}}\left(\left\langle\nabla g, N^{D_{t}}\right\rangle\right)+\bar{\mu}^{\partial D_{t}}(g) \underline{\mu}^{\partial D_{t}}(k) .
\end{aligned}
$$

But $\bar{\mu}^{D_{t}}(\Delta g)=-\bar{\mu}^{\partial D_{t}}\left(\left\langle\nabla g, N^{D_{t}}\right\rangle\right)$ and $F_{k}\left(D_{t}\right)=\mu^{D_{t}}(k)$, so

$$
\begin{equation*}
H\left(D_{t}\right)=\widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g) F_{k}\right) \tag{2.47}
\end{equation*}
$$

which together with (2.45) proves (2.40).
Let us now prove that for any $t \geqslant 0, \mathscr{P}_{t}$ transports $\nu_{0}$ into $\nu_{t}$, where $\left(\mathscr{P}_{t}\right)_{t \geqslant 0}$ is the semi-group introduced after the proof of Proposition 2.10. Consider the map

$$
\begin{equation*}
G(g, k, t)(s)=\nu_{s}\left(\mathscr{P}_{t-s}\left(g F_{k}\right)\right), \quad s \in[0, t] \tag{2.48}
\end{equation*}
$$

We compute

$$
\begin{align*}
G(g, k, t)^{\prime}(s) & =\left(\partial_{s} \nu_{s}\right)\left(\mathscr{P}_{t-s}\left(g F_{k}\right)\right)-\nu_{s}\left(\partial_{t} \mathscr{P}_{t-s}\left(g F_{k}\right)\right) \\
& =\nu_{s}\left(\mathscr{L} \mathscr{P}_{t-s}\left(g F_{k}\right)\right)-\nu_{s}\left(\mathscr{L} \mathscr{P}_{t-s}\left(g F_{k}\right)\right)=0 \tag{2.49}
\end{align*}
$$

where we used Proposition D. 3 in Appendix D to justify the differentiations (as well as the fact that $\mathscr{L} \mathscr{P}_{t-s}\left(g F_{k}\right)=\mathscr{P}_{t-s} \mathscr{L}\left(g F_{k}\right)$ is bounded to be able to use differentiation under the integral $\nu_{s}$ ). So we get $G(g, k, t)(0)=G(g, k, t)(t)$ which rewrites as

$$
\begin{equation*}
\nu_{0} \mathscr{P}_{t}\left(g F_{k}\right)=\nu_{t}\left(g F_{k}\right), \tag{2.50}
\end{equation*}
$$

More generally, by similar arguments, we can replace in this formula $F_{k}$ by any mapping $\mathfrak{F}$ from $\mathfrak{A}$. This in turn implies that $\nu_{0} \mathscr{P}_{t}=\nu_{t}$.

To finish, by iteration, we see that if $X_{0} \sim \bar{\mu}^{D_{0}}$ then $\left(D_{t}\right)_{t \geqslant 0}$ has the same finite time marginals as $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$, proving that $\left(D_{t}\right)$ is a diffusion with generator $\widetilde{\mathscr{L}}$.

## 3. Intertwined dual processes: a generalized Pitman theorem

In this section we will consider the case where $f^{D}$ is the distance to boundary. It is not covered by Section 2 since distance to boundary is not smooth, it is singular on the skeleton of $D$. We will make an approximation of it, and then go to the limit in law.

Let $\widetilde{W}_{t}$ be a real-valued Brownian motion and $\widetilde{D}_{t}$ be the solution of (2.11) started at $\widetilde{D}_{0}$, with driving Brownian motion $\widetilde{W}_{t}$.
Assumption 3.1. Fix $\alpha \in(0,1)$ and $\varepsilon>0$. Let $\widetilde{D}_{0} \in \mathcal{F}^{\alpha, \varepsilon}$. There exists a closed bounded subset $\widetilde{\mathcal{F}}^{\alpha, \varepsilon}$ of $\mathcal{F}^{\alpha, \varepsilon}$ in which the process $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ a.s. takes its values, such that the map $D \mapsto S(D)$ is uniformly Lipschitz from $\widetilde{\mathcal{F}}^{\alpha, \varepsilon}$ with the $C^{2}$ metric to $\mathcal{K}(M)$, the set of compact subsets of $M$ endowed with the Hausdorff metric. Moreover all skeletons $S(D)$ of elements $D \in \widetilde{\mathcal{F}}^{\alpha, \varepsilon}$ have uniformly bounded $(d-1)$-Hausdorff measure, and their regular part have uniformly bounded sectional curvature.

Conjecture 3.2. We conjecture that Assumption 3.1 is always realized, for any $\alpha \in(0,1)$, $\varepsilon>0, \widetilde{D}_{0} \in \mathcal{F}^{\alpha, \varepsilon}$.

All examples together with the study of the motion of the skeleton in Appendix B make us believe that Conjecture 3.2 is true. However a better knowledge of skeletons is necessary to solve it.

Let us begin with some preparatory results. To describe the approximation of $\rho(x, \partial D)$ we are interested in, let us introduce some notations.

- Let $(x, D) \mapsto \ell_{\varepsilon}(x, D):=\left(h_{\varepsilon} \circ \rho_{\partial D}\right)(x)$ where $h_{\varepsilon} \equiv 1$ in $[0, \varepsilon / 2], h_{\varepsilon} \equiv 0$ in $[3 \varepsilon / 4, \infty)$ and $h_{\varepsilon}$ is smooth and nonincreasing in $[0, \infty)$. When $D$ is fixed by the context, we will denote $\ell_{\varepsilon}(x):=\ell_{\varepsilon}(x, D)$.
- For any $\delta \in(0, \varepsilon)$, let $\varphi_{\delta}: \mathbb{R}_{+} \rightarrow \mathbb{R}$ be a nonnegative function with support in $[0, \delta]$, such that the mapping $\mathbb{R}^{d} \ni u \mapsto \varphi_{d}(|u|)$ is smooth and $\int_{\mathbb{R}^{d}} \varphi_{d}(|u|) d u=1$ (in the sequel, $|\cdot|$ will stand for the usual Euclidean norm or for the Riemannian norm on any tangent space of $M$, depending on the context).
- Let $g_{\delta}$ be a smooth, 1-Lipschitz and odd function defined on $\mathbb{R}$, with $g_{\delta}(r)=r$ on $[0, \varepsilon / 4], 0 \leqslant g_{\delta}(r) \leqslant r$ for any $r \geqslant 0$, and $g_{\delta}(r)=c_{\delta} r$ on $[3 \varepsilon / 8, \infty)$, for an appropriate constant $c_{\delta} \leqslant 1$ very close to 1 that will be defined below in (3.2). We write $\rho_{\delta}(x, \partial D):=$ $g_{\delta}(\rho(x, \partial D))$.

The approximation of $\rho(x, \partial D)$ we choose is

$$
\begin{equation*}
f_{\delta}(x, D)=\ell_{\varepsilon}(x, D) \rho_{\delta}(x, \partial D)+\left(1-\ell_{\varepsilon}(x, D)\right) \int_{T_{x} M} \varphi_{\delta}(|v|) \rho_{\delta}\left(\exp _{x}(v), \partial D\right) d v \tag{3.1}
\end{equation*}
$$

(where $d v$ stands for the Lebesgue measure on $T_{x} M$ ).
Since we will need to differentiate $f_{\delta}$, in particular with respect to the first variable (the corresponding gradient will be denoted $\nabla_{1}$ ), it is more convenient to replace the last term of $f_{\delta}(x, D)$ by an integral over the fixed domain $\mathbb{R}^{d}$. More precisely, for any linear isometry (an orthonormal frame) $i_{x}: \mathbb{R}^{d} \rightarrow T_{x} M$, we have

$$
\int_{T_{x} M} \varphi_{\delta}(|v|) \rho_{\delta}\left(\exp _{x}(v), \partial D\right) d v=\int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|) \rho_{\delta}\left(\exp _{x}\left(\imath_{x}(u)\right), \partial D\right) d u
$$

In this purpose, for any $x_{0}, x \in M$, let be given a linear isometry (an orthonormal frame) $\imath_{x_{0}, x}: \mathbb{R}^{d} \rightarrow T_{x} M$, smooth with respect to $x_{0}, x$, and such that $\left.\nabla_{x} l_{x_{0}, x}\right|_{x=x_{0}}=0$. Define

$$
e(\delta):=\sup \left\{\left\|d_{x}\left[\left.\exp _{x}\left(\imath_{x_{0}, x}(u)\right)\right|_{x=x_{0}}\right]\right\| \|, x_{0} \in B(o, R), u \in B(0, \delta) \subset \mathbb{R}^{d}\right\}
$$

where $\left\|\|\cdot\|\right.$ is the operator norm, when $T_{x_{0}} M$ is endowed with its Euclidean structure. Recall that $R=1 / \varepsilon$ and that $\varepsilon$ is fixed as in Assumption 3.1. The previously mentioned constant $c_{\delta}$ is given by

$$
\begin{equation*}
c_{\delta}:=e^{-1}(\delta)\left(1-\delta\left\|\nabla_{1} \ell_{\varepsilon}\right\|_{\infty}\right) \tag{3.2}
\end{equation*}
$$

Notice that $c_{\delta}$ does not depend on $D$ and is as close as we want to 1 .
More precisely, we have
Lemma 3.3. There exists two constants $C_{1}^{\prime}, C_{1}^{\prime \prime}>0$, depending only on $\varepsilon$, such that for $\delta>0$ sufficiently small,

$$
\begin{aligned}
0 \leqslant e(\delta)-1 & \leqslant C_{1}^{\prime} \delta \\
\left|c_{\delta}-1\right| & \leqslant C_{1}^{\prime \prime} \delta
\end{aligned}
$$

Proof. The first inequalities are a consequence of $\left.\nabla_{x} \imath_{x_{0}, x}\right|_{x=x_{0}}=0$ and of the properties of the exponential mapping. The second bound follows, since $\left\|\nabla_{1} \ell_{\varepsilon}\right\|_{\infty}=\left\|h_{\varepsilon}^{\prime}\right\|_{\infty}$ is independent of $D$ (and of order $1 / \varepsilon$ ).

From the second bound, we can and will assume that the function $g_{\delta}$ is furthermore chosen so that $g_{\delta}(r)$ converges uniformly to $r$ on compact sets of $\mathbb{R}_{+}$, as well as the corresponding derivatives up to order 2 as $\delta \searrow 0$. In addition, we choose $\delta>0$ sufficiently small so that the map $(x, y) \mapsto \exp _{x}^{-1}(y)$ is well-defined and smooth in the $\delta$-neighborhood the diagonal of $B(o, R) \times B(o, R)$. Then, for any $x_{0} \in M$, we can rewrite (3.1) under the forms

$$
\begin{align*}
f_{\delta}(x, D) & =\ell_{\varepsilon}(x, D) \rho_{\delta}(x, \partial D)+\left(1-\ell_{\varepsilon}(x, D)\right) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|) \rho_{\delta}\left(\exp _{x_{0}, x}\left(\imath_{x_{0}, x}(u)\right), \partial D\right) d u  \tag{3.3}\\
& =\ell_{\varepsilon}(x, D) \rho_{\delta}(x, \partial D) \\
& +\left(1-\ell_{\varepsilon}(x, D)\right) \int_{M} \varphi_{\delta}\left(\left|\left(\exp _{x}^{\circ} \circ \imath_{x_{0}, x}\right)^{-1}(y)\right|\right) \rho_{\delta}(y, \partial D) J\left(\exp _{x} \circ \imath_{x_{0}, x}\right)^{-1}(y) d y
\end{align*}
$$

where $J\left(\exp _{x} \circ \imath_{x_{0}, x}\right)^{-1}$ is the absolute value of the Jacobian of $\left(\exp _{x} \circ \imath_{x_{0}, x}\right)^{-1}$.
The interest of all these preparations is:

Proposition 3.4. For all $\delta>0$, the function $(x, D) \mapsto f_{\delta}(x, D):=f_{\delta}^{D}(x)$ has the following properties

- $f_{\delta}$ satisfies the conditions of Assumption 2.2;
- there exists $C_{1}>0$ such that $\forall D \in \widetilde{\mathcal{F}}^{\alpha, \varepsilon}$ and $x \in D$, we have

$$
\begin{equation*}
\left|f_{\delta}(x, D)-\rho(x, \partial D)\right| \leqslant C_{1} \delta \tag{3.4}
\end{equation*}
$$

- the gradient and the Hessian of $f_{\delta}$ with respect to the second variable $D$ satisfy $\forall D \in \widetilde{\mathcal{F}}^{\alpha, \varepsilon}, \forall x \in D \backslash S(D)$, for all vector fields $K$ normal to $\partial D$ :

$$
\begin{equation*}
\left\langle\nabla_{2} f_{\delta}(x, D), K\right\rangle \leqslant C_{4}\|K\|_{\infty} \quad \text { and } \quad\left\|\nabla_{2} d_{2} f_{\delta}(x, D)\left(N_{\partial D}, N_{\partial D}\right)\right\| \leqslant C_{4} \tag{3.5}
\end{equation*}
$$

for a $C_{4}$ not depending on $x, D, \delta$.
Proof. We first prove $\left\|\nabla_{1} f_{\delta}(x, D)\right\| \leqslant 1$. Fix $x_{0} \in B(o, R)$ and differentiate with respect to $x$ at $x_{0}$, we have

$$
\begin{align*}
\nabla_{1} f_{\delta}(x, D)= & \ell_{\varepsilon}(x, D) \nabla_{1} \rho_{\delta}(x, \partial D)  \tag{3.6}\\
& +\left(1-\ell_{\varepsilon}(x, D)\right) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|) \nabla_{1} \rho_{\delta}\left(\left(\exp _{x} \circ \imath_{x_{0}, x}\right)(u), \partial D\right) d u \\
& +\nabla_{1} \ell_{\varepsilon}(x, D) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|)\left(\rho_{\delta}(x, \partial D)-\rho_{\delta}\left(\left(\exp _{x} \circ \imath_{x_{0}, x}\right)(u), \partial D\right)\right) d u
\end{align*}
$$

If $\rho(x, \partial D) \leqslant \varepsilon / 2$ then $\ell_{\varepsilon}(x, D)=1, \nabla \ell_{\varepsilon}(x, D)=0$ and

$$
\left\|\nabla_{1} f_{\delta}(x, D)\right\| \leqslant \ell_{\varepsilon}(x, D)\left\|\nabla_{1} \rho_{\delta}(x, \partial D)\right\| \leqslant 1
$$

If $\rho(x, \partial D) \geqslant \varepsilon / 2$ then for $\delta \leqslant \varepsilon / 8$, we have, for $u \in \mathbb{R}^{d}$ with $\left.|u| \leqslant \delta, \rho\left(\exp _{x} \circ \imath_{x_{0}, x}\right)(u), \partial D\right) \geqslant$ $3 \varepsilon / 8$. It follows

$$
\begin{aligned}
\left\|\nabla_{1} f_{\delta}(x, D)\right\| \leqslant & \ell_{\varepsilon}(x) e^{-1}(\delta)\left(1-\delta\left\|\nabla_{1} \ell_{\varepsilon}\right\|_{\infty}\right) \\
& +\left(1-\ell_{\varepsilon}(x)\right) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|) c_{d}\left\|\nabla_{x}\left(\exp _{x} \circ \imath_{x_{0}, x}\right)(u)\right\| d u \\
& +\left\|\nabla_{1} \ell_{\varepsilon}(x)\right\|_{\infty} \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|) \delta d u \\
\leqslant & 1, \quad \text { using } \nabla_{x} \imath_{x_{0}, x}=0
\end{aligned}
$$

It is easily checked that the function $f_{\delta}$ satisfies the other properties of Assumption 2.2. Let us check that it also satisfies (3.4).

We have
$f_{\delta}(x, D)-\rho_{\delta}(x, \partial D)=\left(1-\ell_{\varepsilon}(x, D)\right) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|)\left(\rho_{\delta}\left(\exp _{x}\left(\imath_{x, x}(u)\right), \partial D\right)-\rho_{\delta}(x, \partial D)\right) d u$
which implies

$$
\left|f_{\delta}(x, D)-\rho_{\delta}(x, \partial D)\right| \leqslant \delta
$$

On the other hand

$$
\left|\rho(x, \partial D)-\rho_{\delta}(x, \partial D)\right| \leqslant\left(1-e^{-1}(\delta)\left(1-c_{\delta}\right)\right) 2 R \leqslant C_{1}^{\prime \prime \prime} \delta
$$

for some constant $C_{1}^{\prime \prime \prime}>0$ (depending on $\varepsilon$ ). This yields (3.4) with $C_{1}:=1+C_{1}^{\prime \prime \prime}$.
For proving (3.5), we take a vector field $K(y)=k(y) N(y), y \in \partial D$ and compute

$$
\begin{equation*}
\left\langle\nabla_{2} \rho(x, \partial D), K\right\rangle=\langle-N(P(x)), K(P(x))\rangle=-k(P(x)) \tag{3.8}
\end{equation*}
$$

where $P(x)$ is the projection of $x$ onto $\partial D$, and

$$
\begin{equation*}
\nabla_{2} d_{2} \rho(x, \partial D)\left(N_{\partial D}, N_{\partial D}\right)=0 \tag{3.9}
\end{equation*}
$$

Remarking that $\left\|\nabla_{2} \ell_{\varepsilon}(x, D)\right\|$ is bounded by $\left\|h_{\varepsilon}^{\prime}\right\|_{\infty}$, we get (3.5) via a straightforward computation.

Theorem 3.5. Fix $D_{0}=\tilde{D}_{0} \in \widetilde{\mathcal{F}}^{\alpha, \varepsilon}$ and let $X_{0} \sim \mathscr{U}\left(D_{0}\right)$. Under Assumption 3.1, there exists a pair $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ of $\tau_{\varepsilon}$ intertwined processes in the sense of Definition 1.1, such that the process $\left(D_{t}\right)_{t \geqslant 0}$ satisfies
(3.10)
$d \partial D_{t}(y)$

$$
=N^{D_{t}}(y)\left(\left\langle d X_{t}, N^{D_{t}}\left(X_{t}\right)\right\rangle+\left(\frac{1}{2} h^{D_{t}}(y)-h^{D_{t}}\left(X_{t}\right)\right) d t-2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X)\right)
$$

Here $\theta^{S_{t}}(x)=\pi / 2-\varphi^{S_{t}}(x), \varphi^{S_{t}}(x)$ being the angle between the orthogonal line to $S_{t}$ at $x$ and any of the two minimal geodesics from $\partial D_{t}$ to $x \in S_{t}$. In other words $\theta^{S_{t}}(x)$ is the smallest angle between $S_{t}$ and the geodesics. The process $L^{S_{t}}$ is the local time of $X_{t}$ at $S_{t}:=S\left(D_{t}\right)$ :

$$
\begin{equation*}
L_{t}^{S_{t}}(X)=\lim _{\beta \searrow 0} \frac{1}{2 \beta} \int_{0}^{t} 1_{\left\{X_{s} \in S_{s}^{\beta}\right\}} d s \tag{3.11}
\end{equation*}
$$

$S_{s}^{\beta}$ being the thickening of the regular part of $S_{s}$ in normal direction, of thickness $\beta$ in both directions.

Remark 3.6. - Compared to Section 2 with $f^{D}$ replaced by distance to boundary $\rho_{\partial D}$, we have outside the skeleton $S^{D}$

$$
\begin{equation*}
\nabla \rho_{\partial D}(x)=N^{D}(x) \quad \text { and } \quad \Delta \rho_{\partial D}(x)=-h^{D}(x) \tag{3.12}
\end{equation*}
$$

and we will see that on the moving skeleton $S_{t}=S^{D_{t}}$ :

$$
\begin{equation*}
" \Delta \rho_{\partial D_{t}}\left(X_{t}\right) d t "=-2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X) \tag{3.13}
\end{equation*}
$$

Proof. Under Assumption 3.1, Proposition 3.4 allows us to construct for each $\delta>0$, intertwined processes $\left(X_{t}^{\delta}, D_{t}^{\delta}\right)_{t \geqslant 0}$ started at $\left(X_{0}^{\delta}, D_{0}^{\delta}\right)=\left(X_{0}, D_{0}\right)$, associated with the functions $f_{\delta}^{D}$, stopped at $\tau_{\varepsilon}^{\delta}$, the exit time from $\widetilde{\mathcal{F}}^{\alpha, \varepsilon}$. We have from Equation (2.10)

$$
\begin{equation*}
d \partial D_{t}^{\delta}(y)=N^{D_{t}^{\delta}}(y)\left(d W_{t}^{\delta}+\left(\frac{1}{2} h^{D_{t}^{\delta}}(y)+\Delta f^{D_{t}^{\delta}}\left(X_{t}^{\delta}\right)\right) d t\right) \tag{3.14}
\end{equation*}
$$

for some Brownian motion $W_{t}^{\delta}$. On the other hand, from Proposition 2.11 and (2.1),

$$
\begin{equation*}
\left(\widetilde{D}_{t}^{\delta}\right)_{t \geqslant 0}:=\left(D_{t}^{\delta}\right)_{t \geqslant 0} \tag{3.15}
\end{equation*}
$$

satisfies equation (2.11):

$$
\begin{equation*}
d \partial D_{t}^{\delta}(y)=N^{D_{t}^{\delta}}(y)\left(d \widetilde{W}_{t}^{\delta}+\left(\frac{1}{2} h^{D_{t}^{\delta}}(y)-\frac{\underline{\mu}^{\partial D_{t}^{\delta}}\left(\partial D_{t}^{\delta}\right)}{\mu\left(D_{t}^{\delta}\right)}\right) d t\right) \tag{3.16}
\end{equation*}
$$

where $\widetilde{W}_{t}^{\delta}$ is the $\mathscr{F}_{t}^{D^{\delta}}$-Brownian motion

$$
\begin{equation*}
d \widetilde{W}_{t}^{\delta}=d W_{t}^{\delta}+\Delta f^{D_{t}^{\delta}}\left(X_{t}\right) d t+\frac{\mu^{\partial D_{t}^{\delta}}\left(\partial D_{t}^{\delta}\right)}{\mu\left(D_{t}^{\delta}\right)} d t \tag{3.17}
\end{equation*}
$$

A remarkable fact about all $\left(X_{t}^{\delta}, D_{t}^{\delta}\right)_{t \geqslant 0}$ is that their marginals are constant in law. Notice that also $\left(\left(D_{t}^{\delta}\right)_{t \geqslant 0}, \tau_{\varepsilon}^{\delta}\right)$ is constant in law since $\tau_{\varepsilon}^{\delta}$ is a functional of $\left(D_{t}^{\delta}\right)_{t \geqslant 0}$ independent of $\delta$. As a consequence, the family

$$
\begin{equation*}
\left(\left(X_{t}^{\delta}, D_{t}^{\delta}, W_{t}^{\delta}, \widetilde{W}_{t}^{\delta}, W_{t}^{\delta, m}\right)_{t \geqslant 0}, \tau_{\varepsilon}^{\delta}\right) \tag{3.18}
\end{equation*}
$$

is tight (in (3.18) the Brownian motions $W_{t}^{\delta}$ and $W_{t}^{\delta, m}$ are the ones defined by equation (2.10)). Denote by

$$
\begin{equation*}
\left(\left(X_{t}, D_{t}, W_{t}, \widetilde{W}_{t}, W_{t}^{m}\right)_{t \geqslant 0}, \tau_{\varepsilon}\right) \tag{3.19}
\end{equation*}
$$

a limiting point. Let us prove the intertwining. Using Proposition 2.11, for any smooth functions $g$ and $k$ on $M$,

$$
\begin{aligned}
\mathbb{E}\left[g\left(X_{t}^{\delta}\right) F_{k}\left(D_{t}^{\delta}\right)\right] & =\mathbb{E}\left[\mathbb{E}\left[g\left(X_{t}^{\delta}\right) F_{k}\left(D_{t}^{\delta}\right) \mid \mathscr{F}_{t}^{D^{\delta}}\right]\right] \\
& =\mathbb{E}\left[\mathscr{U}\left(D_{t}^{\delta}\right)(g) F_{k}\left(D_{t}^{\delta}\right)\right] \\
& =\mathbb{E}\left[\frac{F_{g}\left(D_{t}^{\delta}\right)}{F_{1}\left(D_{t}^{\delta}\right)} F_{k}\left(D_{t}^{\delta}\right)\right]
\end{aligned}
$$

and passing to the limit yields the intertwining.
This property of $\left(D_{t}^{\delta}, \widetilde{W}_{t}^{\delta}\right)_{t \geqslant 0}$ being constant in law passes to the limit, and we have

$$
\begin{equation*}
d \partial D_{t}(y)=N^{D_{t}}(y)\left(d \widetilde{W}_{t}+\left(\frac{1}{2} h^{D_{t}}(y)-\frac{\underline{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}\right) d t\right) \tag{3.20}
\end{equation*}
$$

We need to work with real-valued processes: we have from (2.32), for all $\delta>0$,

$$
\begin{equation*}
\int_{0}^{t} \frac{d \mu\left(D_{s}^{\delta}\right)}{\underline{\mu}\left(\partial D_{s}^{\delta}\right)}=-W_{t}^{\delta}-\int_{0}^{t} \Delta_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right) d s \tag{3.21}
\end{equation*}
$$

This together with (3.17) yields

$$
\begin{equation*}
d \partial D_{t}^{\delta}(y)=N^{D_{t}^{\delta}}(y)\left(-\frac{d \mu\left(D_{s}^{\delta}\right)}{\underline{\mu}\left(\partial D_{s}^{\delta}\right)}+\frac{1}{2} h^{D_{t}^{\delta}}(y) d t\right) \tag{3.22}
\end{equation*}
$$

Again by constantness in law:

$$
\begin{equation*}
d \partial D_{t}(y)=N^{D_{t}}(y)\left(-\frac{d \mu\left(D_{s}\right)}{\underline{\mu}\left(\partial D_{s}\right)}+\frac{1}{2} h^{D_{t}}(y) d t\right) \tag{3.23}
\end{equation*}
$$

So to prove our result we only need to prove that

$$
\begin{equation*}
\int_{0}^{t} \frac{d \mu\left(D_{s}\right)}{\underline{\mu}\left(\partial D_{s}\right)}=-W_{t}+\int_{0}^{t} h^{D_{s}}\left(X_{s}\right) d s+\int_{0}^{t} 2 \sin \left(\theta^{S_{s}}\left(X_{s}\right)\right) d L_{s}^{S_{s}}(X) \tag{3.24}
\end{equation*}
$$

and that

$$
\begin{equation*}
W_{t}=\int_{0}^{t}\left\langle N^{D_{s}}\left(X_{s}\right), d X_{s}\right\rangle \tag{3.25}
\end{equation*}
$$

Let us prove (3.25). In all this paragraph we consider $M$ as isometrically embedded in some Euclidean space. In particular we are allowed to integrate vectorial quantities. We use the fact that $d X_{t}^{\delta} \otimes d W_{t}^{\delta}$ converges in law to $d X_{t} \otimes d W_{t}$ (where $\otimes$ stands for bracket of semimartingales). But $d X_{t}^{\delta} \otimes d W_{t}^{\delta}$ is equal to $\nabla_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right) d t$. Then by Lemma H. 1 applied to $\nabla_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right.$ ) (which is uniformly bounded) and $U=\{(x, D), x \notin S(D)\}$ defined in (H.3) we see that the integral of $\nabla_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right) d t$ converges to the one of
$N^{D_{t}}\left(X_{t}\right) d t$. But almost surely $N^{D_{t}}\left(X_{t}\right)$ has norm $1 d t$-a.e., implying that $d W_{t}=$ $\left\langle N^{D_{t}}\left(X_{t}\right), d X_{t}\right\rangle$.

Let us now establish (3.24). It will be a consequence of the convergence of $\left(f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\right)_{t \geqslant 0}$ to $\left(\rho\left(X_{t}, \partial D_{t}\right)_{t \geqslant 0}\right.$.

Write the Itô formula for $f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)$ :

$$
\begin{align*}
d\left(f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\right)= & \left\langle\nabla_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right), d X_{t}^{\delta}\right\rangle+\frac{1}{2} \Delta_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right) d t  \tag{3.26}\\
& +\left\langle\nabla_{2} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right), d \partial D_{t}^{\delta}\right\rangle+\frac{1}{2} \nabla_{2} d_{2} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\left(d \partial D_{t}^{\delta}, d \partial D_{t}^{\delta}\right) d t \\
& +\left\langle\nabla_{21} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right), d \partial D_{t}^{\delta} \otimes d X_{t}^{\delta}\right\rangle
\end{align*}
$$

From Proposition 3.4, possibly by extracting a subsequence,

$$
\begin{equation*}
\left(f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(\rho\left(X_{t}, \partial D_{t}\right)\right)_{t \geqslant 0} . \tag{3.27}
\end{equation*}
$$

From (3.7) we get for $i=1,2$,

$$
\begin{align*}
& \nabla_{i} f_{\delta}(x, D)-\nabla_{i} \rho_{\delta}(x, \partial D) \\
& =-\nabla_{i} \ell_{\varepsilon}(x, D) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|)\left(\rho_{\delta}\left(\exp _{x}\left(\imath_{x, x}(u)\right), \partial D\right)-\rho_{\delta}(x, \partial D)\right) d u  \tag{3.28}\\
& +\left(1-\ell_{\varepsilon}(x, D)\right) \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|)\left(\nabla_{i} \rho_{\delta}\left(\exp _{x}\left(\imath_{x, x}(u)\right), \partial D\right)-\nabla_{i} \rho_{\delta}(x, \partial D)\right) d u
\end{align*}
$$

(where we use the same procedure as in the proof of Proposition 3.4: $\nabla_{1}$ is not applied to the first $x$ of $\imath_{x, x}$ ). From this we see that $\nabla_{1} f_{\delta}(\cdot, D)$ converges, locally uniformly outside $S(D)$, to $\nabla_{1} \rho(\cdot, \partial D)$ with respect to the distance $d_{0}$ of Appendix H. We obtain, with Lemma H.1, possibly by again extracting a subsequence, that

$$
\begin{equation*}
\left(\int_{0}^{t}\left\langle\nabla_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d X_{s}^{\delta}\right\rangle\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(\int_{0}^{t}\left\langle\nabla_{1} \rho\left(X_{s}, \partial D_{s}\right), d X_{s}\right\rangle\right)_{t \geqslant 0} \tag{3.29}
\end{equation*}
$$

More precisely, we have a sequence of martingales converging in law to a martingale $M_{t}$ which is a Brownian motion by Theorem 3 in [21]. For identifying the limiting martingale we use the convergence of $\left\langle\nabla_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d X_{s}^{\delta}\right\rangle \otimes d X_{s}^{\delta}$ to $d M_{s} \otimes d X_{s}$ obtained again by Theorem 3 in [21] (here again we use an isometric embedding of $M$ ). But Lemma H. 1 proves that the limit is equal to $\nabla_{1} \rho\left(X_{s}, \partial D_{s}\right) d s$, yielding (3.29).

Next we prove that

$$
\begin{equation*}
\left(\int_{0}^{t}\left\langle\nabla_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d \partial D_{s}^{\delta}\right\rangle\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(\int_{0}^{t}\left\langle\nabla_{2} \rho\left(X_{s}, \partial D_{s}\right), d \partial D_{s}\right\rangle\right)_{t \geqslant 0} \tag{3.30}
\end{equation*}
$$

The argument is similar except that as we see with (3.14), the drift part of $d \partial D_{s}^{\delta}$ is not well controlled as $X_{t}^{\delta}$ approaches the skeleton. So one cannot proceed exactly the same way. But fortunately, for $x$ outside a $3 \varepsilon / 4$-neighbourhood of $\partial D$ and outside $S(D)$, we have

$$
\begin{align*}
& \left\langle\nabla_{2} f_{\delta}(x, D),\left.N\right|_{\partial D}\right\rangle \\
& =c_{\delta} \int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|)\left\langle-N\left(P\left(\exp _{x}\left(\imath_{x, x}(u)\right)\right), N\left(P\left(\exp _{x}\left(\imath_{x, x}(u)\right)\right)\right\rangle d u=-c_{\delta}\right.\right. \tag{3.31}
\end{align*}
$$

where $c_{\delta}$ is defined in (3.2). This together with (3.22) suggests to write

$$
\begin{aligned}
\int_{0}^{t}\left\langle\nabla_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d \partial D_{s}^{\delta}\right\rangle= & \left(\int_{0}^{t}\left\langle\nabla_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d \partial D_{s}^{\delta}\right\rangle+c_{\delta} \int_{0}^{t} \frac{d \mu\left(D_{s}^{\delta}\right)}{\underline{\mu}\left(\partial D_{s}^{\delta}\right)}\right) \\
& -c_{\delta} \int_{0}^{t} \frac{d \mu\left(D_{s}^{\delta}\right)}{\underline{\mu}\left(\partial D_{s}^{\delta}\right)}
\end{aligned}
$$

The second line clearly converges. The right hand side in the first line can be written

$$
\begin{equation*}
\int_{0}^{t} \tilde{\ell}_{\varepsilon}(x, D)\left\langle\nabla_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right)+\left.c_{\delta} N\right|_{\partial D}, d \partial D_{s}^{\delta}\right\rangle \tag{3.32}
\end{equation*}
$$

with $(x, D) \mapsto \tilde{\ell}_{\varepsilon}(x, D):=\left(\tilde{h}_{\varepsilon} \circ \rho_{\partial D}\right)(x)$ where $\tilde{h}_{\varepsilon} \equiv 1$ in $[0,3 \varepsilon / 4], \tilde{h}_{\varepsilon} \equiv 0$ in $[\varepsilon, \infty)$ and $\tilde{h}_{\varepsilon}$ is smooth and nonincreasing in $[0, \infty)$.

With this last integral we can proceed as for (3.29)
Similarly we obtain the two following convergences for the second derivatives.

$$
\begin{align*}
& \left(\int_{0}^{t} \nabla_{2} d_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right)\left(d \partial D_{s}^{\delta}, d \partial D_{s}^{\delta}\right)\right)_{t \geqslant 0} \\
& \xrightarrow{\mathscr{L}}\left(\int _ { 0 } ^ { t } \nabla _ { 2 } d _ { 2 } \rho ( X _ { s } , \partial D _ { s } ) \left(N\left(P^{\partial D_{s}}\left(X_{s}\right), N\left(P^{\partial D_{s}}\left(X_{s}\right)\right) d s\right)_{t \geqslant 0} \equiv 0\right.\right. \tag{3.33}
\end{align*}
$$

where $P^{\partial D_{s}}\left(X_{s}\right)$ is the orthogonal projection of $X_{s}$ on $\partial D_{s}$ (which is defined $d s$-almost everywhere),

$$
\begin{align*}
& \left.\left(\int_{0}^{t}\left\langle\nabla_{2} d_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d \partial D_{t}^{\delta} \otimes d X_{t}^{\delta}\right\rangle\right)\right)_{t \geqslant 0} \\
& \xrightarrow{\mathscr{L}}\left(\int_{0}^{t}\left\langle\nabla_{2} d_{1} \rho\left(X_{s}, \partial D_{s}\right), d \partial D_{s} \otimes d X_{s}\right\rangle\right)_{t \geqslant 0} \equiv 0 \tag{3.34}
\end{align*}
$$

since $d_{1} \rho\left(X_{s}, \partial D_{s}\right)=-\left\langle N^{D_{s}}\left(X_{s}\right), \cdot\right\rangle$ which implies that the covariant derivative in the second variable with respect to $N^{D_{s}}$ is equal to 0 . On the other hand, by Itô-Tanaka formula, see Proposition F. 1 in Appendix F, using that $\rho(x, \partial D)$ is almost everywhere the minimum of two smooth functions, we have

$$
\begin{align*}
d\left(\rho\left(X_{t}, \partial D_{t}\right)\right)= & \left\langle\nabla_{1} \rho\left(X_{t}, \partial D_{t}\right), d X_{t}\right\rangle-\frac{1}{2} h^{D_{t}}\left(X_{t}\right) d t+\left\langle\nabla_{2} \rho\left(X_{t}, \partial D_{t}\right), d \partial D_{t}\right\rangle  \tag{3.35}\\
& +0+0-\sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X)
\end{align*}
$$

Using (3.26), (3.27), (3.29), (3.30), (3.33), (3.34), (3.35) we obtain that (3.36)

$$
\left(\int_{0}^{t} \Delta_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right) d s\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(\int_{0}^{t}-h^{D_{s}}\left(X_{s}\right) d s-\int_{0}^{t} 2 \sin \left(\theta^{S_{s}}\left(X_{s}\right)\right) d L_{s}^{S_{s}}(X)\right)_{t \geqslant 0}
$$

It remains to pass in the limit as $\delta$ goes to zero in (3.21), to deduce (3.24).
Remark 3.7. From (3.35), it can be deduced that (3.37)

$$
d\left(\rho\left(X_{t}, \partial D_{t}\right)\right)=\frac{1}{2}\left(h^{D_{t}}\left(X_{t}\right)-h^{D_{t}}\left(P^{\partial D_{t}}\left(X_{t}\right)\right)\right) d t+\sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X)
$$

Indeed, (3.25) implies that

$$
\left\langle\nabla_{1} \rho\left(X_{t}, \partial D_{t}\right), d X_{t}\right\rangle=d W_{t}
$$

and due to (3.30), we have

$$
\begin{aligned}
& \left\langle\nabla_{2} \rho\left(X_{t}, \partial D_{t}\right), d \partial D_{t}\right\rangle \\
& \quad=\lim _{\delta \rightarrow 0}\left\langle\nabla_{2} \rho\left(X_{t}^{\delta}, \partial D_{t}^{\delta}\right), d \partial D_{t}^{\delta}\right\rangle \\
& \quad=\lim _{\delta \rightarrow 0}-d W_{t}^{\delta}-\left(\Delta_{1} f_{\delta}\left(P^{\partial D_{t}^{\delta}}\left(X_{t}^{\delta}\right), D_{t}^{\delta}\right)+\frac{1}{2} h^{D_{t}^{\delta}}\left(P^{\partial D_{t}^{\delta}}\left(X_{t}^{\delta}\right)\right)\right) d t
\end{aligned}
$$

where we used (3.21) in conjunction with (3.22).
Taking into account (3.36), we identify the last limit with

$$
-d W_{t}+\left(h^{D_{t}}\left(x_{t}\right)+2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X)-\frac{1}{2} h\left(P^{\partial D_{t}}\left(X_{t}\right)\right)\right) d t
$$

## 4. INTERTWINED DUAL PROCESSES: DECOUPLING AND REFLECTION ON BOUNDARY

In this section we consider another canonical and extremal situation, the case where $f^{D}$ vanishes almost everywhere. More precisely, it is the limiting situation where $f^{D}$ is constant outside a $\varepsilon$-neighbourhood of the boundary. This situation is completely opposite to the one of Section 3 where the coupling is maximal.

Theorem 4.1. There exists a pair $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ of $\tau_{\varepsilon}$-intertwined processes in the sense of Definition 1.1 satisfying

$$
\begin{equation*}
d \partial D_{t}(y) \mathbb{1}_{\left\{X_{t} \notin \partial D_{t}\right\}}=N^{D_{t}}(y)\left(d W_{t}+\frac{1}{2} h^{D_{t}}(y) d t-d L_{t}^{\partial D_{t}}(X)\right) \tag{4.1}
\end{equation*}
$$

where $X_{t}$ is a $M$-valued Brownian motion started at uniform law in $D_{0}, W_{t}$ is a realvalued Brownian motion independent of $X_{t}, L_{t}^{\partial D_{t}}(X)$ is the local time of $X_{t}$ on the moving boundary $\partial D_{t}$.

Remark 4.2. Equation (4.1) can be considered as a limiting case of (2.10). Here Assumption 3.1 is not needed since the morphological skeleton of $D$ does not play a role, and the map $D \mapsto \partial D$ is already sufficiently regular.

Proof. The proof is quite similar to the one of Theorem 3.5, but with another family of functions $f_{\delta}^{D}$, namely $f_{\delta}^{D}:=h_{\delta} \circ \rho_{\partial D}$ where $h_{\delta}$ is defined in the proof of Proposition 2.4: $h_{\delta}$ is a smooth nondecreasing function from $[0, \infty)$ to $\mathbb{R}_{+}$such that $h_{\delta}(r)=r$ for $r \in$ $[0, \delta / 2], h_{\delta}(r)=(3 / 4) \delta$ for $r \geqslant \delta$ and $\left\|h_{\delta}^{\prime}\right\|_{\infty} \leqslant 1$. But here, as $\varepsilon$ is fixed, we will let $\delta \searrow 0$. Again we construct for each $\delta>0$, an intertwined processes $\left(X_{t}^{\delta}, D_{t}^{\delta}\right)_{t \geqslant 0}$ stopped at $\tau_{\varepsilon}^{\delta}$. Again all $\left(X_{t}^{\delta}, D_{t}^{\delta}\right)_{t \geqslant 0}$ are tight, and a limiting process $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ stopped at $\tau_{\varepsilon}$ provides an intertwining. The proof of (4.1) goes along the same lines as the one of (3.10).

We end this section with another canonical construction, where the functions $f_{\delta}^{D}$ approximate $-\rho_{\partial D}$.

Theorem 4.3. Under assumption 3.4, there exists an intertwining $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ stopped at $\tau_{\varepsilon}$, satisfying

$$
\begin{align*}
d \partial D_{t}(y)= & N^{D_{t}}(y)\left(-\left\langle d X_{t}, N^{D_{t}}\left(X_{t}\right)\right\rangle+\left(\frac{1}{2} h^{D_{t}}(y)+h^{D_{t}}\left(X_{t}\right)\right) d t\right.  \tag{4.2}\\
& \left.+2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X)-2 d L_{t}^{\partial D_{t}}(X)\right)
\end{align*}
$$

Proof. It is completely similar to the ones of Theorems 3.5 and 4.1.

## 5. Some fundamental examples

5.1. Real Brownian motion and three-dimensional Bessel process. We come back to the case where $M=\mathbb{R}$. Assume that the Brownian motion $X$ starts from 0 (to respect rigorously the above framework, $X$ should start from the uniform distribution on $D_{0}:=$ $[-\epsilon, \epsilon]$ and next we should let $\epsilon$ go to $0_{+}$). Due to the invariance by symmetry of (3.10), for any $t>0, D_{t}$ remains a symmetric interval, let us write it $\left[-R_{t}, R_{t}\right]$. In this simple setting, we have $N^{D_{t}}(\cdot)=-\operatorname{sign}(\cdot)$ on $\mathbb{R} \backslash\{0\}, h^{D_{t}}=0$ and $S_{t}=\{0\}$, for any $t>0$. Thus (3.10) writes

$$
\begin{equation*}
d R_{t}=\operatorname{sign}\left(X_{t}\right) d X_{t}+2 d L_{t} \tag{5.1}
\end{equation*}
$$

where $L:=\left(L_{t}\right)_{t \geqslant 0}$ is the local time of $X$ at 0 . Namely we get that

$$
\begin{aligned}
\forall t \geqslant 0, \quad R_{t} & =\int_{0}^{t} \operatorname{sign}\left(X_{s}\right) d X_{s}+2 L_{t} \\
& =\left|X_{t}\right|+L_{t}
\end{aligned}
$$

by Tanaka's formula. It is well-known that $R:=\left(R_{t}\right)_{t \geqslant 0}$ is a Bessel process of dimension 3 (cf. e.g. Corollary 3.8 of Chapter 6 of Revuz and Yor [17]). In particular, we get that with the notation introduced in (A.4),

$$
\forall t \geqslant 0, \quad \rho_{\partial D_{t}}^{+}\left(X_{t}\right)=\min \left(X_{t}+R_{t}, R_{t}-X_{t}\right)
$$

But except at time $t=0$, this quantity is always positive: a.s. $X_{t}$ never touch the boundary of $D_{t}$ for $t>0$. Indeed, if for some $t>0$ we have $\left|X_{t}\right|=R_{t}$, we deduce that $L_{t}=0$, namely a contradiction, since $X_{0}=0$.

In particular, we see that the intertwining coupling we have constructed is different from the one proposed by Pitman [16], which is a.s. touching (the upper) boundary repeatedly. Instead we end up with the intertwining dual constructed in [14] via stochastic flows. It is mentioned there how to deduce the classical Pitman's dual, via Lévy's theorem.

Here is an alternative approach. While Equation (5.1) is obtained from approximating $x \mapsto|r-x|$ outside an $\varepsilon$-neighbourhood of 0 when $D=[-r, r]$ by smooth functions $f^{D}$ satisfying Assumption 2.2, we are able to recover Pitman theorem by rather approximating $x \mapsto-x$ in $D=[-r, r]$ outside the only $\varepsilon$-neighbourhood of $-r$. In the limit of (2.10) as $\epsilon$ goes to zero, on the one hand we have

$$
\begin{equation*}
\mathbb{1}_{\left\{X_{t} \neq R_{t}\right\}} d R_{t}=d X_{t} \tag{5.2}
\end{equation*}
$$

on the other hand we have $X_{t}+R_{t} \geqslant 0$, so that $X_{t}+R_{t}$ is the solution to the Skorohod problem associated to $2 X_{t}$. We get

$$
\begin{equation*}
R_{t}+X_{t}=2 X_{t}-2 \min _{0 \leqslant s \leqslant t} X_{s} \tag{5.3}
\end{equation*}
$$

which is equivalent to

$$
\begin{equation*}
R_{t}=X_{t}-2 \min _{0 \leqslant s \leqslant t} X_{s} \tag{5.4}
\end{equation*}
$$

The answer to the question: what would be a symmetric construction with local time at the two ends of $\left[-R_{t}, R_{t}\right]$ is given by Theorem 4.3. We obtained intertwined processes with

$$
\begin{equation*}
R_{t}=-\int_{0}^{t} \operatorname{sign}\left(X_{s}\right) d X_{s}-2 L_{t}^{0}(X)+2 L_{t}^{0}(R-X)+2 L_{t}^{0}(R+X) \tag{5.5}
\end{equation*}
$$

5.2. Brownian motion and disks in rotationally symmetric manifolds. This is the most simple example since the skeleton is never hit by the Brownian motion. Consider a complete $d$-dimensional manifold with $d \geqslant 2$, rotationally symmetric around a point $o \in M$. Denote by $(r, \Theta)$ polar coordinates with $r(x)=\rho(o, x)$ and

$$
\begin{equation*}
d s^{2}=d r^{2}+f^{2}(r) d \Theta^{2} \tag{5.6}
\end{equation*}
$$

the metric in polar coordinates. Then the radial Laplacian is

$$
\begin{equation*}
\Delta_{r}=\frac{\partial^{2}}{(\partial r)^{2}}+b(r) \frac{\partial}{\partial r} \quad \text { with } \quad b=(d-1)(\ln f)^{\prime} \tag{5.7}
\end{equation*}
$$

We will investigate set-valued processes $D_{t}=B\left(o, R_{t}\right)$ where $B(o, r)$ is the open geodesic ball centered at $o$, with radius $r$. The skeleton of $B\left(o, R_{t}\right)$ is the point $o$.

Let $X_{t}$ be a Brownian motion in $M$ satisfying $X_{0} \sim \mathscr{U}\left(D_{0}\right)$ for some $D_{0}=B\left(o, r_{0}\right)$. Denote by $\rho_{t}:=r\left(X_{t}\right)$ the radial part of $X_{t}$. Then

$$
\begin{equation*}
d \rho_{t}=d \beta_{t}+\frac{1}{2} b\left(\rho_{t}\right) d t, \quad \rho_{0} \sim \mathscr{U}^{f}\left(\left(0, r_{0}\right)\right) \tag{5.8}
\end{equation*}
$$

where $\left(\beta_{t}\right)_{t \geqslant 0}$ is a real Brownian motion and

$$
\begin{equation*}
\mathscr{U}^{f}(d r):=\frac{f(r)}{\int_{0}^{r_{0}} f(s) d s} d r \tag{5.9}
\end{equation*}
$$

The evolution equation (3.10) for $D_{t}$ shows by symmetry that for all $t \geqslant 0, D_{t}=$ $B\left(0, R_{t}\right)$ for some real-valued process $R_{t}$. Moreover it writes

$$
\begin{align*}
d \rho_{t} & =d \beta_{t}+\frac{1}{2} b\left(\rho_{t}\right) d t \\
d R_{t} & =d \beta_{t}+\left[-\frac{1}{2} b\left(R_{t}\right)+b\left(\rho_{t}\right)\right] d t \tag{5.10}
\end{align*}
$$

Proposition 5.1. The system of equations (5.10) has a solution up to explosion time of $R_{t}$

$$
\begin{equation*}
\tau^{D}:=\inf \left\{t \geqslant 0, R_{t} \notin(0, \infty)\right\} \tag{5.11}
\end{equation*}
$$

which satisfies for all $t<\tau^{D}$,

$$
\begin{equation*}
0<\rho_{t}<R_{t} \tag{5.12}
\end{equation*}
$$

The corresponding set-valued process $D_{t}=B\left(o, R_{t}\right)$ is solution to equation (3.10), and in particular, for all $\mathscr{F}^{D}$-stopping time $\tau$,

$$
\begin{equation*}
\mathscr{L}\left(X_{\tau} \mid \mathscr{F}_{\tau}^{D}\right)=\mathscr{U}\left(D_{\tau}\right) \quad \text { as well as } \quad \mathscr{L}\left(\rho_{\tau} \mid \mathscr{F}_{\tau}^{D}\right)=\mathscr{U}^{f}\left(\left(0, R_{\tau}\right)\right) \tag{5.13}
\end{equation*}
$$

Proof. We only have to check (5.12). By (5.10),

$$
\begin{equation*}
d\left(R_{t}-\rho_{t}\right)=\frac{1}{2}\left[b\left(\rho_{t}\right)-b\left(R_{t}\right)\right] d t \tag{5.14}
\end{equation*}
$$

which vanishes on $\left\{R_{t}=\rho_{t}\right\}$, and since $b$ is smooth, if $\rho_{0}<R_{0}$, then $\rho_{t}<\mathbb{R}_{t}$ for all times.

### 5.3. Brownian motion and annulus in 2-dimensional rotationally symmetric mani-

 folds. Let $M$ be a complete 2 -dimensional Riemannian manifold, rotationally symmetric around a point $o \in M$. Denote by $(r, \theta)$ polar coordinates with $r(x)=\rho(o, x)$ and$$
\begin{equation*}
d s^{2}=d r^{2}+f^{2}(r) d \theta^{2} \tag{5.15}
\end{equation*}
$$

the metric in polar coordinates. Then the radial Laplacian is

$$
\begin{equation*}
\Delta_{r}=\frac{\partial^{2}}{(\partial r)^{2}}+b(r) \frac{\partial}{\partial r} \quad \text { with } \quad b=(\ln f)^{\prime} \tag{5.16}
\end{equation*}
$$

If $0 \leqslant r^{-} \leqslant r^{+}$, let

$$
\begin{equation*}
A\left(r^{-}, r^{+}\right):=\left\{x \in M, r^{-}<r(x)<r^{+}\right\} \quad \text { if } \quad r^{-}<r^{+}, \quad A\left(r^{-}, r^{-}\right):=\varnothing \tag{5.17}
\end{equation*}
$$

the open annulus delimited by the radius $r^{-}$and $r^{+}$.
In the following we will investigate set-valued processes $D_{t}=A\left(R_{t}^{-}, R_{t}^{+}\right)$. The skeleton of $A\left(R_{t}^{-}, R_{t}^{+}\right)$is the circle

$$
\begin{equation*}
S_{t}=C\left(o, R_{t}^{0}\right) \quad \text { with } \quad R_{t}^{0}:=\frac{1}{2}\left(R_{t}^{-}+R_{t}^{+}\right) \tag{5.18}
\end{equation*}
$$

Let $X_{t}$ be a Brownian motion in $M$ satisfying $X_{0} \sim \mathscr{U}\left(D_{0}\right)$ for some $D_{0}=A\left(r_{0}^{-}, r_{0}^{+}\right)$. Denote by $\rho_{t}:=r\left(X_{t}\right)$ the radial part of $X_{t}$. Then

$$
\begin{equation*}
d \rho_{t}=d \beta_{t}+\frac{1}{2} b\left(\rho_{t}\right) d t, \quad \rho_{0} \sim \mathscr{U}^{f}\left(\left(r_{0}^{-}, r_{0}^{+}\right)\right) \tag{5.19}
\end{equation*}
$$

where $\beta_{t}$ is a real Brownian motion and

$$
\begin{equation*}
\mathscr{U}^{f}\left(\left(r_{0}^{-}, r_{0}^{+}\right)\right)(d r):=\frac{f(r)}{\int_{r_{0}^{-}}^{r_{0}^{+}} f(s) d s} d r . \tag{5.20}
\end{equation*}
$$

The evolution equation (3.10) for $D_{t}$ shows by symmetry that for all $t \geqslant 0, D_{t}=$ $A\left(R_{t}^{-}, R_{t}^{+}\right)$for some real-valued processes $R_{t}^{-} \leqslant R_{t}^{+}$. Moreover it writes

$$
\begin{align*}
d \rho_{t} & =\operatorname{sign}\left(\rho_{t}-R_{t}^{0}\right) d W_{t}+\frac{1}{2} b\left(\rho_{t}\right) d t \\
d R_{t}^{+} & =d W_{t}+\left[-\frac{1}{2} b\left(R_{t}^{+}\right)+\operatorname{sign}\left(\rho_{t}-R_{t}^{0}\right) b\left(\rho_{t}\right)\right] d t+2 L_{t}^{R_{t}^{0}}(\rho) \\
d R_{t}^{-} & =-d W_{t}+\left[-\frac{1}{2} b\left(R_{t}^{-}\right)-\operatorname{sign}\left(\rho_{t}-R_{t}^{0}\right) b\left(\rho_{t}\right)\right] d t-2 L_{t}^{R_{t}^{0}}(\rho)  \tag{5.21}\\
R_{t}^{0} & =\frac{1}{2}\left(R_{t}^{-}+R_{t}^{+}\right)
\end{align*}
$$

and these equations imply

$$
\begin{equation*}
d R_{t}^{0}=-\frac{1}{4}\left[b\left(R_{t}^{+}\right)+b\left(R_{t}^{-}\right)\right] d t \tag{5.22}
\end{equation*}
$$

Proposition 5.2. The system of equations (5.21) has a solution up to explosion time

$$
\begin{equation*}
\tau^{D}:=\inf \left\{t \geqslant 0,\left(R_{t}^{-}, R_{t}^{+}\right) \notin(0, \infty)^{2}\right\} \tag{5.23}
\end{equation*}
$$

which satisfies for all $t<\tau^{D}$,

$$
\begin{equation*}
R_{t}^{-} \leqslant \rho_{t} \leqslant R_{t}^{+} \tag{5.24}
\end{equation*}
$$

The corresponding set-valued process $D_{t}=A\left(R_{t}^{-}, R_{t}^{+}\right)$is solution to equation (3.10), and in particular, for all $\mathscr{F}^{D}$-stopping time $\tau$,

$$
\begin{equation*}
\mathscr{L}\left(X_{\tau} \mid \mathscr{F}_{\tau}^{D}\right)=\mathscr{U}\left(D_{\tau}\right) \quad \text { as well as } \quad \mathscr{L}\left(\rho_{\tau} \mid \mathscr{F}_{\tau}^{D}\right)=\mathscr{U}^{f}\left(\left(R_{\tau}^{-}, R_{\tau}^{+}\right)\right) . \tag{5.25}
\end{equation*}
$$

Proof. Fix $\varepsilon>0$ and $\alpha \in(0,1)$. We will first solve the system of equations until the exit time $\tau_{\varepsilon}$ and then let $\varepsilon \searrow 0$. Let us construct functions $f_{\delta}^{D}(x)$ which satisfies equation (3.1). It will be easier here because there is no need of functions $\ell_{\varepsilon}$ and $g_{\delta}$.

For $\delta \in(0, \varepsilon)$, let $\varphi_{\delta}: \mathbb{R} \rightarrow \mathbb{R}$ be the function with support equal to [ $-\delta / 2, \delta / 2$ ], satisfying for $-\delta / 2<r<\delta / 2$ :

$$
\begin{equation*}
\varphi_{\delta}(r):=\frac{1}{c(\delta)} \exp \left(-\frac{1}{\left(\frac{\delta}{2}\right)^{2}-r^{2}}\right) \quad \text { with } \quad c(\delta):=\int_{-\delta / 2}^{\delta / 2} \exp \left(-\frac{1}{\left(\frac{\delta}{2}\right)^{2}-s^{2}}\right) d s \tag{5.26}
\end{equation*}
$$

and let

$$
\begin{align*}
\operatorname{sign}_{\delta}: \mathbb{R} & \rightarrow \mathbb{R} \\
r & \mapsto-1+2 \int_{-\infty}^{r} \varphi_{\delta}(s) d s \tag{5.27}
\end{align*}
$$

The functions $\varphi_{\delta}$ and $\operatorname{sign}_{\delta}$ are both smooth and Lipschitz, and they respectively approximate $\delta_{0}$ and sign. For $0<r^{-}<r^{+}$satisfying $r^{+}-r^{-} \geqslant 2 \varepsilon$, defining $r^{0}:=\frac{1}{2}\left(r^{-}+r^{+}\right)$, for $x \in A\left(r^{-}, r^{+}\right)$let
$f^{A\left(r^{-}, r^{+}\right)}(x)=f\left(x, r^{-}, r^{+}\right)=g(r(x)) \quad$ with $\quad g(r)=g\left(r, r^{-}, r^{+}\right)=\int_{r^{-}}^{r}-\operatorname{sign}_{\delta}\left(s-r^{0}\right) d s$.
Clearly $f\left(x, r^{-}, r^{+}\right)$is 1-Lipschitz in the first variable. A computation shows that

$$
\begin{equation*}
\partial_{r^{+}} g\left(r, r^{-}, r^{+}\right)=\int_{-\varepsilon}^{r^{0}} \varphi_{\delta}(v) d v \quad \text { and } \quad \partial_{r^{-}} g\left(r, r^{-}, r^{+}\right)=-\int_{r-r^{0}}^{\varepsilon} \varphi_{\delta}(v) d v \tag{5.29}
\end{equation*}
$$

showing that $g$ and $f$ are 1-Lipschitz. Then the vector $N:=N_{\partial A\left(r^{-}, r^{+}\right)}$is equal to $-\mathbb{1}_{\left\{r(x)=r^{+}\right\}} \partial_{r^{+}}+\mathbb{1}_{\left\{r(x)=r^{-}\right\}} \partial_{r^{-}}$so that

$$
\begin{equation*}
\langle\nabla f, N\rangle \equiv 1 \quad \text { and } \quad \nabla d f(N, N) \equiv 0 \tag{5.30}
\end{equation*}
$$

This yields an elementary proof of the properties of Proposition 3.4. We can use Theorem 3.5 to solve equation (5.21) until the stopping time $\tau_{\varepsilon}$.

We are left to prove that $\tau_{\varepsilon} \nearrow \tau^{D}$ a.s. as $\varepsilon \searrow 0$. This is a direct consequence of the fact that the volume of $A\left(R_{t}^{-}, R_{t}^{+}\right)$is a time changed Bessel process of dimension 3 (by [6] Theorem 5), proving that $A\left(R_{t}^{-}, R_{t}^{+}\right)$cannot collapse onto its skeleton.

Remark 5.3. After the hitting time of 0 by $R_{t}^{-}$, the processes can continue to evolve under the regime of Section 5.2.

We recover from Proposition 5.2 a result from [13] stating that $\left(\left[R_{t}^{-}, R_{t}^{+}\right]\right)_{t \geqslant 0}$ is an intertwining dual process for the real diffusion $\left(\rho_{t}\right)_{t \geqslant 0}$. In particular, we deduce that if $\left(\rho_{t}\right)_{t \geqslant 0}$ is positive recurrent and if $+\infty$ is an entrance boundary, then $\left(\left[R_{t}^{-}, R_{t}^{+}\right]\right)_{t \geqslant 0}$ reaches $[0,+\infty]$ in finite time and this finite time is a strong stationary time for $\left(\rho_{t}\right)_{t \geqslant 0}$, see [13] for more details.
5.4. Brownian motion and symmetric convex sets in $\mathbb{R}^{2}$. In this section we take $M=$ $\mathbb{R}^{2}$ endowed with the Euclidean metric. Consider a smooth strictly convex bounded set $D_{0} \subset M$ with smooth boundary, symmetric with respect to the horizontal and vertical axes. Also assume that its skeleton is an horizontal interval $S_{0}^{\prime}=\left[-x_{0}, x_{0}\right] \times\{0\}$. An example of such a set is the interior of an ellipse, the skeleton being the interval between the two foci. Assume that $X_{t}$ is a Brownian motion in $\mathbb{R}^{2}$ satisfying $X_{0} \sim \mathscr{U}\left(D_{0}\right)$. Let us investigate the evolution of $\left(X_{t}, D_{t}\right)$. Notice that it is the first example where we really have to deal with infinite dimensional processes. By conservation of the convexity by the normal and mean curvature flows, $D_{t}$ will stay convex. It will also stay symmetric. Propositions 5.4, 5.5 and 5.7 below will be proved in a forthcoming paper:
Proposition 5.4. The skeleton of $D_{t}$ always takes the form $\left[-x_{t}, x_{t}\right] \times\{0\}$.
Proof. This will be proved in [1]
Denote by $(\imath, \jmath)$ the canonical basis of $\mathbb{R}^{2}$, and $X_{t}=\left(X_{t}^{(1)}, X_{t}^{(2)}\right)$. In this notation, the vector $N^{D_{t}}\left(X_{t}\right)$ of Equation (3.10) writes

$$
\begin{equation*}
N^{D_{t}}\left(X_{t}\right)=-\operatorname{sign}\left(X_{t}^{(1)}\right) \cos \left(\theta^{S_{t}}\left(X_{t}\right)\right) \imath-\operatorname{sign}\left(X_{t}^{(2)}\right) \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) \jmath \tag{5.31}
\end{equation*}
$$

where $\theta^{S_{t}}(x)$ is naturally extended to $D_{t}$ by being constant on lines normal to the boundary (see [1]). Notice that $x \mapsto \theta^{S_{t}}(x)$ is locally Lipschitz on $D_{t}$ and is equal to 0 on $D_{t} \cap$ $\left(\left(-\infty,-x_{t}\right] \times\{0\} \cup\left[x_{t}, \infty\right) \times\{0\}\right)$. Also notice that the function $h^{D_{t}}$ is locally Lipschitz on $D_{t} \backslash\left\{\left(-x_{t}, 0\right),\left(x_{t}, 0\right)\right\}$. With these notations, equation (3.10) writes

$$
\begin{align*}
d \partial D_{t}(y)= & -N^{D_{t}}(y)\left(\operatorname{sign}\left(X_{t}^{(1)}\right) \cos \left(\theta^{S_{t}}\left(X_{t}\right)\right) d X_{t}^{(1)}+\sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) \operatorname{sign}\left(X_{t}^{(2)}\right) d X_{t}^{(2)}\right.  \tag{5.32}\\
& \left.+\left(\frac{1}{2} h^{D_{t}}(y)-h^{D_{t}}\left(X_{t}\right)\right) d t-2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}\left(X^{(2)}\right)\right)
\end{align*}
$$

Let us investigate the motion of the skeleton $\tilde{S}_{t}$ of the solution $\tilde{D}_{t}$ of equation (2.11) (garanteed by Theorem 2.5).
Proposition 5.5. The process $\left(\tilde{D}_{t}\right)_{t \geqslant 0}$ takes its values in a closed subset $\tilde{\mathcal{F}}^{\alpha, \varepsilon}$ of $\mathcal{F}^{\alpha, \varepsilon}$, symmetric with respect to the horizontal and vertical axes, such that on $\widetilde{\mathcal{F}}^{\alpha, \varepsilon}$, the map $\left.D \mapsto h^{D}\right|_{\partial D}$ is continuous from $\tilde{\mathcal{F}}^{\alpha, \varepsilon}$ (with the $C^{2+\alpha}$ metric) to $C^{2}(\partial D)$. Its skeleton $\tilde{S}_{t}$ satisfies $\tilde{S}_{t}=\left[-\tilde{x}_{t}, \tilde{x}_{t}\right] \times\{0\}$ for some process $\tilde{x}_{t}$.
Proof. This will be proved in [1]
In the next result we prove that the skeleton has finite variation and is monotonly decreasing.
Proposition 5.6. The right endpoint $\left(\tilde{x}_{t}, 0\right)$ of the skeleton $\tilde{S}_{t}$ satisfies

$$
\begin{equation*}
\frac{d \tilde{x}_{t}}{d t}=\frac{\rho^{2}\left(\left(\tilde{x}_{t}, 0\right), \tilde{y}_{t}\right)}{2}\left(h^{\tilde{D}_{t}}\right)^{\prime \prime}\left(\tilde{y}_{t}\right) \tag{5.33}
\end{equation*}
$$

$\tilde{y}_{t}$ being the point of $\partial \tilde{D}_{t}$ in the horizontal line with the greatest abscissa, and the second derivative being calculated with curvilign coordinates on $\partial \tilde{D}_{t}$. Notice that $\left(h^{\tilde{D}_{t}}\right)^{\prime \prime}\left(\tilde{y}_{t}\right) \leqslant 0$, proving that the process $S\left(\tilde{D}_{t}\right)$ is monotonly decreasing.

Proof. Let us investigate the motion of a point in $\tilde{S}_{t}$ close to $\left(\tilde{x}_{t}, 0\right)$. This point has two closest points in $\partial \tilde{D}_{t}$, which we call $\widetilde{y}_{1, t}$ and $\widetilde{y}_{2, t}$, the first one having positive second coordinate. We will use Theorem B. 1 and (B.28). Call $\hat{x}_{t}$ the point in the skeleton corresponding to $\widetilde{y}_{1, t}$ and $\widetilde{y}_{2, t}$. We have $N_{1}\left(\hat{x}_{t}\right)=-\cos \theta\left(\hat{x}_{t}\right) \imath-\sin \theta\left(\hat{x}_{t}\right) \jmath, N_{2}\left(\hat{x}_{t}\right)=$ $-\cos \theta\left(\hat{x}_{t}\right) \imath+\sin \theta\left(\hat{x}_{t}\right) \jmath, N_{1}^{S}\left(\hat{x}_{t}\right)=-\jmath$. Denote $T\left(\widetilde{y}_{1, t}\right)$ the tangent vector to $\partial \tilde{D}_{t}$ at $\widetilde{y}_{1, t}$, corresponding to increasing of $\theta: T\left(\widetilde{y}_{1, t}\right)=-\sin \theta\left(\hat{x}_{t}\right) \imath+\cos \theta\left(\hat{x}_{t}\right) \jmath$. Write $h^{\prime}\left(\widetilde{y}_{1, t}\right)$ the curvilign derivative of $h\left(\widetilde{y}_{1, t}\right)$ in the direction of $T\left(\widetilde{y}_{1, t}\right)$. Then the vector $J_{1}^{\perp}(1)$ of (B.28) is equal to $-\frac{1}{2} \rho_{S}\left(\widetilde{y}_{1, t}\right) h^{\prime}\left(\widetilde{y}_{1, t}\right) T\left(\widetilde{y}_{1, t}\right)$. So we get from (B.28):

$$
\begin{align*}
\frac{d}{d t} \hat{x}_{t} & =\frac{1}{2} \rho_{S}\left(\widetilde{y}_{1, t}\right) h^{\prime}\left(\widetilde{y}_{1, t}\right)\left(\sin \theta\left(\hat{x}_{t}\right)+\frac{\cos ^{2} \theta\left(\hat{x}_{t}\right)}{\sin \theta\left(\hat{x}_{t}\right)}\right) \imath \\
& =\frac{\rho_{S}\left(\widetilde{y}_{1, t}\right) h^{\prime}\left(\widetilde{y}_{1, t}\right)}{2 \sin \theta\left(\hat{x}_{t}\right)} \imath  \tag{5.34}\\
& =\frac{\rho_{S}^{2}\left(\widetilde{y}_{1, t}\right) h^{\prime}\left(\widetilde{y}_{1, t}\right)}{2 \widetilde{y}_{1, t}^{(2)}} \imath \quad \text { with } \quad \widetilde{y}_{1, t}=\left(\widetilde{y}_{1, t}^{(1)}, \widetilde{y}_{1, t}^{(2)}\right) .
\end{align*}
$$

In the limit, as $\tilde{y}_{1, t}^{(2)}$ goes to zero, we obtain the motion of $\tilde{x}_{t}$ and using the symmetry of the convex set, we have $h^{\prime}\left(\widetilde{y}_{t}\right)=0$ so that we can replace $\frac{h^{\prime}\left(\widetilde{y}_{1, t}\right)}{\widetilde{y}_{1, t}^{(2)}}$ by $h^{\prime \prime}\left(\widetilde{y}_{t}\right)$. This yields (5.33).

In particular a Brownian motion $X_{t}$ will never meet the ends of $\tilde{S}_{t}$.
A solution to (5.32) can be found with the help of Theorem 3.5. The family of functions $f_{\delta}(x, D)$ defined in (3.1) takes the form:

$$
\begin{align*}
f_{\delta}(x, D) & =\ell_{\varepsilon}(x) \rho_{\delta}(x, \partial D)+\left(1-\ell_{\varepsilon}(x)\right) \int_{\mathbb{R}^{2}} \varphi_{\delta}(|x-y|) \rho_{\delta}(y, \partial D) d y \\
& =\ell_{\varepsilon}(x) \rho_{\delta}(x, \partial D)+\left(1-\ell_{\varepsilon}(x)\right) \int_{\mathbb{R}^{2}} \varphi_{\delta}(|y|) \rho_{\delta}(x-y, \partial D) d y \tag{5.35}
\end{align*}
$$

Proposition 5.7. Equation (5.32) provides an intertwining with infinite lifetime
Proof. This will be proved in [1]

## Appendix A. An integration by parts on domains with boundary

Let $M$ be a $d$-dimensional Riemannian manifold and $D \subset M$ a compact and connected domain with smooth boundary $\partial D$. For $y \in \partial D$, let $N(y)$ be the inward normal vector. Denote by $S^{\prime}$ the inward (morphological) skeleton of $D: S^{\prime}$ is the set of points in $D$ such that the distance to $\partial D$ is not smooth with non vanishing gradient around them. Denote

$$
\begin{equation*}
\tau(y)=\inf \left\{t>0, \exp _{y}(t N(y)) \in S^{\prime}\right\} \tag{A.1}
\end{equation*}
$$

Let $S$ be the set of regular points of $S^{\prime}$, which we can describe as follows: if $x \in S$, then there exists a unique couple $\left(y_{1}, y_{2}\right)$ of distinct points from $\partial D$ such that

$$
\begin{equation*}
x=\exp _{y_{1}}\left(\tau\left(y_{1}\right) N\left(y_{1}\right)\right)=\exp _{y_{2}}\left(\tau\left(y_{2}\right) N\left(y_{2}\right)\right) . \tag{A.2}
\end{equation*}
$$

We have $\tau\left(y_{1}\right)=\tau\left(y_{2}\right)$, and for $i=1,2$, the differential at $\left(\tau\left(y_{i}\right), y_{i}\right)$ of the map $\mathbb{R}_{+} \times$ $\partial D \ni(t, y) \mapsto \exp _{y}(t N(y))$ is nondegenerate. The set $S$ is a codimension 1 submanifold of $M$ and $S^{\prime} \backslash S$ has Hausdorff dimension smaller than or equal to $d-2$. It is the union of the focal set which is the set of points $x=\exp _{y}(\tau(y) N(y))$ such that $\left(t, y^{\prime}\right) \mapsto \exp _{y^{\prime}}\left(t N\left(y^{\prime}\right)\right)$ is degenerate at $(\tau(y), y)$, and the union of the sets defined like $S$ but withstrictly more than two points $y_{1}, y_{2}, y_{3}, \ldots$ For $r \geqslant 0$, let

$$
\begin{equation*}
D(r)=\left\{z \in D \backslash S^{\prime}, \rho_{\partial D}(z) \geqslant r\right\} . \tag{A.3}
\end{equation*}
$$

where $\rho$ is the Riemannian distance. The set $D(r)$ is a (possibly empty) manifold with smooth boundary $\partial D(r)$ on which one can define an inward normal $N(y)$ and an orientation by parallel transporting oriented basis of $\partial D$ along normal geodesics. So we have for all $y \in D \backslash S^{\prime}: N(y)=\nabla \rho_{\partial D}(y)$.

We will also need the sets $D(r)$ for all $r \in \mathbb{R}$. We will let for $r<0$

$$
\begin{equation*}
D(r)=\left\{z \in M, \rho_{\partial D}^{+}(z) \geqslant r\right\} \tag{A.4}
\end{equation*}
$$

where $\rho_{\partial D}^{+}$is the signed distance to $\partial D$, positive inside $D$, negative outside $D$.
Define for $s, t \in \mathbb{R}$

$$
\begin{align*}
\psi(s, t): \partial D(s) & \rightarrow \partial D(t) \\
y & \mapsto \exp _{y}((t-s) N(y)) \tag{A.5}
\end{align*}
$$

and $\psi(t)=\psi(0, t)$. We will indifferentely write $\psi(t)(x)=\psi(t, x)$. The function $\psi(s, t)$ is not defined for all points of $\partial D(s)$ because we ask $\psi(s, t)(y) \in \partial D(t)$, nor is $N(\cdot)$. However for $|s|$ and $|t|$ small it is a map, defined for all $y \in \partial D(s)$, and is is also a diffeomorphism with inverse $\psi(t, s)$.

We have for $0 \leqslant s \leqslant t, \psi(t)=\psi(s, t) \circ \psi(s)$, which implies

$$
\begin{equation*}
\operatorname{det} T \psi(t)=\operatorname{det} T \psi(s, t) \times \operatorname{det} T \psi(s) \tag{A.6}
\end{equation*}
$$

Notice that thanks to the orientation of the sets $\partial D(r)$ we get an orientation of $D \backslash S^{\prime}$ by adding $N$ as first vector to oriented basis, consequently $\operatorname{det} T \psi$ is well defined and always positive. It is well-known that

$$
\begin{equation*}
\left.\frac{d}{d t}\right|_{t=s} \operatorname{det} T \psi(s, t)(y)=-h(y) \tag{A.7}
\end{equation*}
$$

where $h(y)$ is the inward mean curvature of $\partial D(s)$ (the minus sign of the r.h.s. of (A.7) insures that $h$ is non-negative on $\partial D(s)$ when $D(s)$ is convex). This together with (A.6) yields

$$
\begin{equation*}
\left.\frac{d}{d t}\right|_{t=s} \operatorname{det} T \psi(t)(y)=-h(\psi(s)(y)) \operatorname{det} T \psi(s)(y) \tag{A.8}
\end{equation*}
$$

and consequently, using $\psi(0)=\mathrm{id}$ and $\operatorname{det} T \psi(0) \equiv 1$,

$$
\begin{equation*}
\operatorname{det} T \psi(t)(y)=\exp \left(\int_{0}^{t}-h(\psi(s)(y)) d s\right) \tag{A.9}
\end{equation*}
$$

Denote by $\mu$ the volume measure of $D$ and by $\mu$ the volume measures of the manifolds $\partial D(s)$ and of $S$. Then

$$
\begin{equation*}
\mu(D)=\int_{0}^{\infty} \underline{\mu}(\partial D(r)) d r . \tag{A.10}
\end{equation*}
$$

But for $r \geqslant 0$

$$
\begin{equation*}
\underline{\mu}(\partial D(r))=\int_{\partial D} \operatorname{det} T \psi(r)(y) \underline{\mu}(d y) \tag{A.11}
\end{equation*}
$$

with convention $\operatorname{det} T \psi(r)(y)=0$ if $r \geqslant \tau(y)$. We get

$$
\begin{equation*}
\underline{\mu}(\partial D(r))=\int_{\partial D} \exp \left(-\int_{0}^{r} h(\psi(s)(y)) d s\right) 1_{\{r<\tau(y)\}} \underline{\mu}(d y) \tag{A.12}
\end{equation*}
$$

which yields with (A.10)

$$
\begin{equation*}
\mu(D)=\int_{\partial D}\left(\int_{0}^{\tau(y)} \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right) d r\right) \underline{\mu}(d y) \tag{A.13}
\end{equation*}
$$

More generally, for a measurable function $g: D \rightarrow \mathbb{R}$ bounded below,

$$
\begin{equation*}
\int_{D} g d \mu=\int_{\partial D}\left(\int_{0}^{\tau(y)} g(\psi(r, y)) \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right) d r\right) \underline{\mu}(d y) \tag{A.14}
\end{equation*}
$$

Applying this formula to the function $g h$ which we assume to be bounded below or integrable, we get by integration by parts

$$
\begin{aligned}
\int_{D} g h d \mu & =\int_{\partial D}\left(\int_{0}^{\tau(y)}-g(\psi(r, y)) \frac{d}{d r} \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right) d r\right) \underline{\mu}(d y) \\
& =\int_{\partial D}\left[-g(\psi(r, y)) \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right)\right]_{0}^{\tau(y)} \underline{\mu}(d y) \\
& +\int_{\partial D}\left(\int_{0}^{\tau(y)}\langle d g, N\rangle(\psi(r, y)) \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right) d r\right) \underline{\mu}(d y) \\
& =\int_{\partial D} g(y) \underline{\mu}(d y)-\int_{\partial D} g(\psi(\tau(y), y)) e^{-\int_{0}^{\tau(y)} h(\psi(u, y)) d u} \underline{\mu}(d y) \\
& +\int_{D}\langle d g, N\rangle d \mu
\end{aligned}
$$

Define the map

$$
\begin{align*}
\varphi: \partial D & \rightarrow S^{\prime} \\
y & \mapsto \psi(\tau(y), y) \tag{A.15}
\end{align*}
$$

For $z=\psi\left(\tau\left(y_{i}\right), y_{i}\right) \in S(i=1,2)$ define $\theta(z) \in(0, \pi / 2]$ the angle between $N\left(\psi\left(\tau\left(y_{i}\right)-, y_{i}\right)\right)$ and $S$. In the sequel we assume that $\theta(z) \neq \pi / 2$ (the case $\theta(z)=\pi / 2$ is simpler to deal with and Proposition A. 1 is always valid). Notice that this angle does not depend on $i$, this is a consequence of $z \in S$ staying at the same distance to $y_{1}$ and $y_{2}$ by infinitesimal variation. For later use, let also $\theta(z)=0$ when $z \in S^{\prime} \backslash S$. Let us prove that for $\left.z=\psi\left(\tau\left(y_{i}\right), y_{i}\right)\right) \in S$,
(A.16) $\quad \operatorname{det} T \psi\left(\tau\left(y_{i}\right), y_{i}\right)=\sin \theta\left(\varphi\left(y_{i}\right)\right) \operatorname{det} T \varphi\left(y_{i}\right), \quad i=1,2$.

Set $y=y_{1}$. Let $e_{1}=N(y), e_{1}^{S}=N(\psi(\tau(y)-, y)), N^{S}(z)$ the normal to $S$ at $z$ such that $\left\langle N^{S}(z), e_{1}^{S}\right\rangle>0$, let $e^{\prime \prime}=\left(e_{3}, \ldots, e_{d}\right)$ be a family of orthonormal normalized vectors in $T_{y} \partial D$ such that letting $e_{2}=\frac{\nabla \tau(y)}{\|\nabla \tau(y)\|}$ (we have $\nabla \tau(y) \neq 0$, since $\theta(z) \neq \pi / 2$ ), $e^{\prime}:=$ $\left(e_{2}, e^{\prime \prime}\right)$ is an orthonormal basis of $T_{y} \partial D$, let $\left(e^{S}\right)^{\prime \prime}=\left(e_{3}^{S}, \ldots, e_{d}^{S}\right)$ be an orthonormal basis of $T_{y} \varphi\left(\operatorname{Vect}\left(e^{\prime \prime}\right)\right)$, let $e_{2}^{S}$ such that $\left(e^{S}\right)^{\prime}:=\left(e_{2}^{S}, \ldots, e_{d}^{S}\right)$ is an orthonormal basis of $T_{z} S$. Finally let $e_{2}^{\theta} \in T_{z} M$ be such that $\left\langle e_{2}^{\theta}, N(z)\right\rangle<0\left(e_{2}^{\theta}\right.$ and $N^{S}(z)$ are not orthogonal, since $\theta(z) \neq \pi / 2)$ and $\left(e_{1}^{S}, e_{2}^{\theta},\left(e^{S}\right)^{\prime \prime}\right)$ is an orthonormal basis of $T_{z} M$. Figure 1 shows the configuration of $e_{1}^{S}, N^{S}(z), e_{2}^{S}$ and $e_{2}^{\theta}$ on an example of dimension 2. In the sequel we
will denote for instance $T \varphi\left(e^{\prime}\right)=\left(\begin{array}{c}T \varphi\left(e_{2}\right) \\ \vdots \\ T \varphi\left(e_{d}\right)\end{array}\right)$, so that $\left\langle T \varphi\left(e^{\prime}\right),\left(e^{S}\right)^{\prime}\right\rangle$ will be the matrix of all scalar products. We have

$$
\begin{aligned}
& \left\langle T \varphi\left(e^{\prime}\right),\left(e^{S}\right)^{\prime}\right\rangle \\
& =\left\langle d \tau, e^{\prime}\right\rangle\left\langle\partial_{t} \psi(\tau(y), y),\left(e^{S}\right)^{\prime}\right\rangle+\left\langle T \psi\left(e^{\prime}\right),\left(e^{S}\right)^{\prime}\right\rangle \\
& =\left(\begin{array}{ll}
\left\langle d \tau, e_{2}\right\rangle\left\langle\partial_{t} \psi, e_{2}^{S}\right\rangle+\left\langle T \psi\left(e_{2}\right), e_{2}^{S}\right\rangle & \left\langle T \psi\left(e_{2}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle \\
\left\langle d \tau, e^{\prime \prime}\right\rangle\left\langle\partial_{t} \psi, e_{2}^{S}\right\rangle+\left\langle T \psi\left(e^{\prime \prime}\right), e_{2}^{S}\right\rangle & \left\langle T \psi\left(e^{\prime \prime}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle
\end{array}\right) .
\end{aligned}
$$

Let us simplify and make more explicit this expression. We have $\left\langle d \tau, e^{\prime \prime}\right\rangle=0$. Also $e_{2}^{\theta} \perp\left(e^{S}\right)^{\prime \prime}$ and $e_{2}^{S} \perp\left(e^{S}\right)^{\prime \prime}$ so $e_{2}^{S} \in \operatorname{Vect}\left(e_{1}^{S}, e_{2}^{\theta}\right)$ and more precisely

$$
\begin{equation*}
e_{2}^{S}=\cos (\theta(z)) e_{1}^{S}+\sin (\theta(z)) e_{2}^{\theta} \tag{A.17}
\end{equation*}
$$

On the other hand $T \psi\left(e^{\prime}\right) \perp e_{1}^{S}$ which implies

$$
\begin{equation*}
\left\langle T \psi\left(e^{\prime}\right), e_{2}^{S}\right\rangle=\sin (\theta(z))\left\langle T \psi\left(e^{\prime}\right), e_{2}^{\theta}\right\rangle \tag{A.18}
\end{equation*}
$$

Also $\left\langle\partial_{t} \psi, e_{2}^{S}\right\rangle=\cos (\theta(z))$. We arrive at

$$
\begin{align*}
& \operatorname{det}\left\langle T \varphi\left(e^{\prime}\right),\left(e^{S}\right)^{\prime}\right\rangle \\
& =\sin \theta(z) \operatorname{det}\left(\begin{array}{cc}
\left\langle T \psi\left(e_{2}\right), e_{2}^{\theta}\right\rangle & \left\langle T \psi\left(e^{\prime \prime}\right), e_{2}^{\theta}\right\rangle \\
\left\langle T \psi\left(e_{2}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle & \left\langle T \psi\left(e^{\prime \prime}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle
\end{array}\right)  \tag{A.19}\\
& +\cos \theta(z) \operatorname{det}\left(\begin{array}{cc}
\left\langle d \tau, e_{2}\right\rangle & 0 \\
\left\langle T \psi\left(e_{2}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle & \left\langle T \psi\left(e^{\prime \prime}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle
\end{array}\right) \\
& =\sin \theta(z) \operatorname{det} T \psi+\cos \theta(z)\left\langle d \tau, e_{2}\right\rangle \operatorname{det}\left\langle T \psi\left(e^{\prime \prime}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle .
\end{align*}
$$

For the last equation we used the fact that $\operatorname{det} T \psi=\operatorname{det}\left\langle T \psi\left(e^{\prime}\right),\left(e_{2}^{\theta},\left(e^{S}\right)^{\prime \prime}\right)\right\rangle$, since $e^{\prime}$ and $\left(e_{2}^{\theta},\left(e^{S}\right)^{\prime \prime}\right)$ are orthonormal bases. Note that by definition, $\left\langle T \psi\left(e^{\prime \prime}\right), e_{2}^{\theta}\right\rangle=0$, so we also get $\operatorname{det} T \psi=\operatorname{det}\left\langle T \psi\left(e^{\prime \prime}\right),\left(e^{S}\right)^{\prime \prime}\right\rangle \times\left\langle T \psi\left(e_{2}\right), e_{2}^{\theta}\right\rangle$. On the other hand, we have

$$
\begin{equation*}
\left\langle d \tau, e_{2}\right\rangle=\left\langle T \psi\left(e_{2}\right), e_{2}^{\theta}\right\rangle \cot \theta(z) \tag{A.20}
\end{equation*}
$$

Indeed, note that

$$
\begin{aligned}
0 & =\left\langle T \varphi\left(e_{2}\right), N^{S}\right\rangle \\
& =\left\langle d \tau, e_{2}\right\rangle\left\langle e_{1}^{S}, N^{S}\right\rangle+\left\langle T \psi\left(e_{2}\right), N^{S}\right\rangle \\
& =\left\langle d \tau, e_{2}\right\rangle \sin (\theta(z))-\cos (\theta(z))\left\langle T \psi\left(e_{2}\right), e_{2}^{\theta}\right\rangle
\end{aligned}
$$

where the last term is obtained by taking into account that $T \psi\left(e_{2}\right)$ is parallel to $e_{2}^{\theta}$. This is the change of length of the geodesic needed to stay in $S$. We obtain

$$
\begin{aligned}
\operatorname{det} T \varphi & =\sin \theta(z) \operatorname{det} T \psi+\cos \theta(z) \cot \theta(z) \operatorname{det} T \psi \\
& =\frac{\sin ^{2} \theta(z)+\cos ^{2} \theta(z)}{\sin \theta(z)} \operatorname{det} T \psi
\end{aligned}
$$

This yields (A.16).
We arrived at

$$
\begin{align*}
\int_{D} g h d \mu & =\int_{\partial D} g(y) \underline{\mu}(d y)-\int_{\partial D} g(\psi(\tau(y), y)) \operatorname{det} T \psi(\tau(y), y) \underline{\mu}(d y)  \tag{A.21}\\
& +\int_{D}\langle d g, N\rangle d \mu
\end{align*}
$$



Figure 1. The vectors $e_{1}^{S}, N^{S}(z), e_{2}^{S}$ and $e_{2}^{\theta}$
this yields with (A.16)

$$
\begin{align*}
\int_{D} g h d \mu & =\int_{\partial D} g(y) \underline{\mu}(d y)-\int_{\partial D} g(\varphi(y)) \sin \theta(\varphi(y)) \operatorname{det} T \varphi(y) \underline{\mu}(d y)  \tag{A.22}\\
& +\int_{D}\langle d g, N\rangle d \mu
\end{align*}
$$

Using the change of variable $y \mapsto \varphi(y)$ and the fact that all $z \in S$ is equal to $\varphi\left(y_{i}\right)$, $i=1,2$, we obtain the key formula

Proposition A.1. With the above notations, for any smooth function $g$ defined on $D$ such that gh is integrable or bounded below, we have:

$$
\begin{equation*}
\int_{D} g h d \mu=\int_{\partial D} g(y) \underline{\mu}(d y)-2 \int_{S} g(z) \sin \theta(z) \underline{\mu}(d z)+\int_{D}\langle d g, N\rangle d \mu \tag{A.23}
\end{equation*}
$$

## Appendix B. Moving Sets

In this section we describe how to move a domain with smooth boundary by deformation of its boundary. We will investigate the deformation of its skeleton The deformation we will consider will have a general absolutely continuous finite variation part, together with a very specific martingale part and singular finite variation part. First we introduce some notation.

For a domain $D$ with smooth boundary $\partial D, s \in \mathbb{R}$, define

$$
\begin{align*}
\psi^{D}(s)=\psi^{D}(0, s): \partial D & \rightarrow \partial D(s) \\
y & \mapsto \psi^{D}(s)(y)=\psi^{D}(s, y)=\exp _{y}\left(s N^{D}(y)\right) \tag{B.1}
\end{align*}
$$

Here $N^{D}=N$ is the inward normal defined in Section A. Consider a moving domain $t \mapsto D_{t}$. Be careful not to confound $D(t)$ with $D_{t}$, since in general they are quite different subsets. We first assume that the deformation is sufficiently regular so that for all $0 \leqslant s \leqslant$ $t$, we can write $D_{t}$ as

$$
\begin{equation*}
D_{t}=\left\{\psi^{D_{s}}\left(\left[Z_{t}^{D_{s}}(y), \tau_{D_{s}}(y)\right], y\right), \quad y \in \partial D_{s}\right\} \tag{B.2}
\end{equation*}
$$

In particular, we must have $S_{s}^{\prime} \subset D_{t}$. Notice that in the special case where the real valued function $t \mapsto Z_{t}^{D_{s}}(y)$ does not depend on $y$, for any $0 \leqslant s \leqslant t$, then we have

$$
\begin{equation*}
D_{t}=D_{s}\left(Z_{t}^{D_{s}}\right)=D_{0}\left(Z_{t}^{D_{0}}\right), \quad Z_{t}^{D_{0}}=Z_{t}^{D_{s}}+Z_{s}^{D_{0}} \tag{B.3}
\end{equation*}
$$

where $D(r)$ is defined in (A.3), replacing distance to $\partial D$ by signed distance with positive sign inside $D$ and negative sign outside. In this situation, the skeleton is not moving, at least as long as $\partial D_{t}$ remains smooth (i.e. until $\partial D_{t}$ hits $S_{0}^{\prime}$ or is too far outside $D_{0}$ ), and $t \mapsto Z_{t}^{D_{0}}$ can be allowed to be a semimartingale with singular continuous drift.

When $t \mapsto Z_{t}^{D_{s}}(y)$ depends on $y$ the situation is a little bit more complicated. Starting from $(t, y) \mapsto Z_{t}^{D_{0}}(y)$ which is assumed to be defined on $[0, \varepsilon) \times \partial D_{0}$, the sets $D_{t}$ are defined for $0 \leqslant t<\varepsilon$, as well as the $Z_{t}^{D_{s}}(y), 0 \leqslant s \leqslant t, y \in D_{s}$. In fact, if $(y, t) \mapsto Z_{t}^{D_{0}}(y)$ is $C^{1}$, then one can reconstruct all $Z_{t}^{D_{s}}(y)$ with the only knowledge of $\dot{Z}_{t}^{D_{t}}(z), z \in \partial D_{t}$. Let us do it for $s=0$ : the map $(t, y) \mapsto \psi^{D_{0}}(t, y)$ from $(-\alpha, \alpha) \times \partial D_{0}$ to $M$ is a diffeomorphism on its range, for $\alpha>0$ sufficiently small. Let us denote $z \mapsto\left(\tau_{0}(z), \varphi_{0}(z)\right)$ its inverse. Then a variation $z+N^{D_{t}}(z) d Z_{t}^{D_{t}}$ corresponds to a variation $\left(\tau_{0}(z), \varphi_{0}(z)\right)+\left(d \tau_{0}, T \varphi_{0}\right) N^{D_{t}}(z) d Z_{t}^{D_{t}}$ of the coordinates in $(-\alpha, \alpha) \times \partial D_{0}$. But this is not convenient at all, since it is not intrinsic. Moreover, when passing to stochastic processes and Stratonovich equations, it will involve second derivatives of $z \mapsto\left(\tau_{0}(z), \varphi_{0}(z)\right)$. So we prefer to leave the reference to $D_{0}$ and to always stay at the level of the moving $D_{t}$.

For all $y \in \partial D_{0}$ we define a stochastic process $t \mapsto Y_{t}(y)$ representing the motion of $D_{t}$ satisfying $Y_{0}(y)=y$ and the Itô equation in manifold with respect to the Levi Civita connection $\nabla$
(B.4)

$$
d Y_{t}(y)=d^{\nabla} Y_{t}(y)=\partial_{1} \psi^{D_{t}}\left(\cdot, Y_{t}(y)\right)\left(d Z_{t}^{D_{t}}\left(Y_{t}(y)\right)\right)=N^{D_{t}}\left(Y_{t}(y)\right) d Z_{t}^{D_{t}}\left(Y_{t}(y)\right)
$$

Recall that fomally $d^{\nabla} Y_{t}(y)$ is a vector which writes in local coordinates $\left(y^{1}, \ldots, y^{d}\right)$ with the Christoffel symbols $\Gamma_{j, k}^{i}$ :

$$
\begin{equation*}
d^{\nabla} Y_{t}(y)=\left(d Y_{t}^{i}(y)+\frac{1}{2} \Gamma_{j, k}^{i}\left(Y_{t}(y)\right) d\left\langle Y_{t}^{j}(y), Y_{t}^{k}(y)\right\rangle\right) D_{i}\left(Y_{t}(y)\right) \tag{B.5}
\end{equation*}
$$

where $D_{i}\left(Y_{t}(y)\right)$ is the vector $\frac{\partial}{\partial y^{i}}$ taken at point $Y_{t}(y)$. We will always assume that the martingale part $d m_{t}$ of $d Z_{t}^{D_{t}}(y)$ does not depend on $y$. In this situation, the Itô equation is equivalent to the Stratonovich one: indeed, using (B.3) the Itô to Stratonovich convertion term is

$$
\frac{1}{2} \nabla_{N^{D_{t}}\left(Y_{t}(y)\right) d m_{t}} N^{D_{t}}(\cdot) d m_{t}=\frac{1}{2} \nabla_{N^{D_{t}}\left(Y_{t}(y)\right)} N^{D_{t}}(\cdot) d\langle m, m\rangle_{t}=0
$$

since $N^{D_{t}}\left(Y_{t}(y)\right)$ is the speed at time $a=0$ of the geodesic $a \mapsto \psi^{D_{t}}(a)\left(Y_{t}(y)\right)$.
More precisely, we will let $d Z_{t}^{D_{t}}(y)$ be of the form

$$
\begin{equation*}
d Z_{t}^{D_{t}}(y)=H^{D_{t}}\left(Y_{t}(y)\right) d t+d z_{t} \tag{B.6}
\end{equation*}
$$

where $H^{D_{t}}$ is a smooth function on $\partial D_{t}$ (which later on will be chosen to be $h^{D_{t}} / 2$, where $h^{D_{t}}$ is the mean curvature of $\left.\partial D_{t}\right)$ and $\left(z_{t}\right)_{t \geqslant 0}$ is a real valued continuous semimartingale. We assume that Equation (B.4) has a strong solution up to some positive stopping time. Moreover, since $d Y_{t}(y)$ represents the motion of $\partial D_{t}$ and for small time the map $y^{\prime} \mapsto$ $Y_{t}\left(y^{\prime}\right)$ is a diffeomorphism from $\partial D_{0}$ to $\partial D_{t}$, writing $Y_{t}\left(y^{\prime}\right)=y$, equation (B.4) rewrites as

$$
\begin{equation*}
d \partial D_{t}(y):=d Y_{t}\left(y^{\prime}\right)=N^{D_{t}}(y)\left(H^{D_{t}}(y) d t+d z_{t}\right) \tag{B.7}
\end{equation*}
$$

Let us now investigate the motion of the skeleton $S_{t}$ under this motion of $D_{t}$. First we remark that by local inversion theorem, at regular points of the skeleton, the variation in Stratonovich sense is linear and the sum of all variations of the concerned point at the boundary. As we already remarked, the motion $d z_{t}$ does not change $S_{t}$, so this together with the linearity just mentioned implies that we have a finite variation of the skeleton.

Recall the situation of (A.2) in Section A. We consider a domain $D, x \in S, y_{1}, y_{2}$ the two elements of $\partial D$ such that $\exp _{y_{1}}\left(\tau\left(y_{1}\right) N\left(y_{1}\right)\right)=\exp _{y_{2}}\left(\tau\left(y_{2}\right) N\left(y_{2}\right)\right)$, with $\tau\left(y_{1}\right)=$ $\tau\left(y_{2}\right)$. For $i=1,2$, we will consider a variation of the minimal geodesic from $y_{i}$ to $x$, represented by a Jacobi field $J_{i}$ satisfying $J_{i}(0) \in T_{y_{i}} M, J_{1}(1)=J_{2}(1) \in T_{x} M$,

$$
\begin{equation*}
J_{i}(0)=\lambda_{i} N\left(y_{i}\right)+J_{i}^{\perp}(0), \quad J_{i}^{\prime}(0)=\lambda_{i}^{\prime} N\left(y_{i}\right)+\left(J_{i}^{\perp}\right)^{\prime}(0) \tag{B.8}
\end{equation*}
$$

with $J_{i}^{\perp}$ orthogonal to $N\left(y_{i}\right)$. The motion of $S$ corresponding to the motion of $y_{1}$ and $y_{2}$ will be represented by $J_{1}(1)$. Since $S$ has a boundary, the observation of the orthogonal part to $S$ of $J_{1}(1)$ is not sufficient.

Let $\gamma_{i}$ be the projection on $M$ of $J_{i}$. It is the geodesic in time 1 from $y_{i}$ to $x$ (as usual in the computations of Jacobi fields, the speed is not normalized). Denote $N_{i}(x)=$ $\dot{\gamma}_{i}(1) /\left\|\dot{\gamma}_{i}(1)\right\|$. Recall that the angle between $N_{i}(x)$ and $T_{x} S$ is $\theta(x) \in(0, \pi / 2]$. We will also let

$$
\begin{equation*}
N_{1}^{S}(x)=\frac{1}{2 \sin \theta(x)}\left(N_{1}(x)-N_{2}(x)\right) \tag{B.9}
\end{equation*}
$$

Figure 2 shows the configuration of the points $x, y_{1}, y_{2}$ and the vectors $N_{1}(x), N_{2}(x)$, $N_{1}^{S}(x)$. The vector $N_{1}^{S}(x)$ is is the normal vector to $S$ at point $x$, in the same side as $N_{1}(x)$. We will consider variations of geodesics with same final value:

$$
\begin{equation*}
J_{1}(1)=J_{2}(1)=\lambda N_{1}^{S}(x)+J_{1}^{T}(1) \tag{B.10}
\end{equation*}
$$



Figure 2. The points $x, y_{1}, y_{2}$ and the vectors $N_{1}(x), N_{2}(x), N_{1}^{S}(x)$
for some $\lambda \in \mathbb{R}$, where $J_{1}^{T}(1) \in T_{x} S$. Writing $\lambda N_{1}^{S}(x)=\frac{\lambda}{2 \sin \theta(x)}\left(N_{1}(x)-N_{2}(x)\right)$ we have
(B.11)

$$
\begin{aligned}
\left\langle J_{1}(1), N_{1}(x)\right\rangle & =\frac{\lambda}{2 \sin \theta(x)}\left(1-\cos (2 \theta(x))+\left\langle J_{1}^{T}(1), N_{1}(x)\right\rangle\right. \\
& =\lambda \sin \theta(x)+\left\langle J_{1}^{T}(1), N_{1}(x)\right\rangle
\end{aligned}
$$

and

$$
\begin{align*}
\left\langle J_{1}(1), N_{2}(x)\right\rangle & =-\frac{\lambda}{2 \sin \theta(x)}\left(1-\cos (2 \theta(x))+\left\langle J_{1}^{T}(1), N_{2}(x)\right\rangle\right.  \tag{B.12}\\
& =-\lambda \sin \theta(x)+\left\langle J_{1}^{T}(1), N_{2}(x)\right\rangle
\end{align*}
$$

On the other hand we require that the variation of length of the two geodesics are the same. This writes as

$$
\begin{equation*}
\left\langle J_{1}(1), N_{1}(x)\right\rangle-\left\langle J_{1}(0), N\left(y_{1}\right)\right\rangle=\left\langle J_{2}(1), N_{2}(x)\right\rangle-\left\langle J_{2}(0), N\left(y_{2}\right)\right\rangle \tag{B.13}
\end{equation*}
$$

or
(B.14) $\quad \lambda \sin \theta(x)+\left\langle J_{1}^{T}(1), N_{1}(x)\right\rangle-\lambda_{1}=-\lambda \sin \theta(x)+\left\langle J_{1}^{T}(1), N_{2}(x)\right\rangle-\lambda_{2}$, which finally, with $\left\langle J_{1}^{T}(1), N_{1}(x)-N_{2}(x)\right\rangle=0$, yields $\lambda=\frac{\lambda_{1}-\lambda_{2}}{2 \sin \theta(x)}$, so the normal variation of $S$ is given by

$$
\begin{equation*}
\left\langle J_{1}(1), N_{1}^{S}(x)\right\rangle N_{1}^{S}(x)=\frac{\lambda_{1}-\lambda_{2}}{2 \sin \theta(x)} N_{1}^{S}(x) . \tag{B.15}
\end{equation*}
$$

Next we will compute the tangential displacement $J^{T}(1)$ of $x$ in $S$. As we will see later, we will only need a Jacobi field $J_{1}$ such that $J_{1}^{\perp}(0)$ and $\left(J_{1}^{\perp}\right)^{\prime}(0)$ are known and

$$
\begin{equation*}
J_{1}(0)=\lambda_{1} N\left(y_{1}\right), \text { i.e. } \quad J_{1}^{\perp}(0)=0 . \tag{B.16}
\end{equation*}
$$

So we know $J_{1}^{\perp}(1)$ : and

$$
\begin{equation*}
J_{1}^{\perp}(1)=J\left(1,0,\left(J_{1}^{\perp}\right)^{\prime}(0)\right) \tag{B.17}
\end{equation*}
$$

where $J(1, u, v)$ is the value at time 1 of the Jacobi field $J$ with $J(0)=u$ and $J^{\prime}(0)=v$. From

$$
\begin{align*}
& J_{1}(1)=J_{1}^{T}(1)+\left\langle J_{1}(1), N_{1}^{S}(x)\right\rangle N_{1}^{S}(x)  \tag{B.18}\\
& J_{1}(1)=J_{1}^{\perp}(1)+\left\langle J_{1}(1), N_{1}(x)\right\rangle N_{1}(x)
\end{align*}
$$

we get

$$
\begin{equation*}
J_{1}^{T}(1)=J_{1}^{\perp}(1)+\left\langle J_{1}(1), N_{1}(x)\right\rangle N_{1}(x)-\left\langle J_{1}(1), N_{1}^{S}(x)\right\rangle N_{1}^{S}(x) \tag{B.19}
\end{equation*}
$$

On the other hand we have

$$
\begin{align*}
& \left\langle J_{1}(1), N_{2}(x)\right\rangle=\left\langle J_{1}^{\perp}(1), N_{2}(x)\right\rangle+\left\langle J_{1}(1), N_{1}(x)\right\rangle\left\langle N_{1}(x), N_{2}(x)\right\rangle \\
& \left\langle J_{1}(1), N_{2}(x)\right\rangle=\left\langle J_{1}(1), N_{1}(x)\right\rangle-\left(\lambda_{1}-\lambda_{2}\right) \tag{B.20}
\end{align*}
$$

where the second equation is a direct consequence of (B.15). Substracting the second equation to the first one yields

$$
\begin{equation*}
(1-\cos (2 \theta(x)))\left\langle J_{1}(1), N_{1}(x)\right\rangle=\left\langle J_{1}^{\perp}(1), N_{2}(x)\right\rangle+\lambda_{1}-\lambda_{2} \tag{B.21}
\end{equation*}
$$

Replacing $\left\langle J_{1}(1), N_{1}(x)\right\rangle$ in (B.19) and after simplification, using (B.9) and (B.15), we finally obtain the horizontal displacement
(B.22)
$\left(J_{1}^{T}\right)(1)=J_{1}^{\perp}(1)+\frac{1}{4 \sin ^{2} \theta(x)}\left(2\left\langle J_{1}^{\perp}(1), N_{2}(x)\right\rangle N_{1}(x)+\left(\lambda_{1}-\lambda_{2}\right)\left(N_{1}(x)+N_{2}(x)\right)\right)$.
We are now in position to write the motion of the skeleton $S_{t}$ when the motion of the boundary is given by (B.7). For $x \in S_{t}$ with corresponding points $y_{1}$ and $y_{2}$ in $\partial D_{t}$,

$$
\begin{equation*}
d S_{t}^{\perp}(x)=\frac{1}{2 \sin \theta^{S_{t}}(x)}\left(H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)\right) N_{1}^{S_{t}}(x) d t \tag{B.23}
\end{equation*}
$$

which has finite variation. Observe that, as already mentioned, the term $d z_{t}$ disappears.
Here we wrote $d S_{t}^{\perp}(x)$ for the normal variation of the regular skeleton. But as we already remarked, since $S_{t}$ is not a closed manifold, it can expand via the motion of its boundary. So we have to investigate the horizontal motion $d S^{T}(x)$.

Notice that $\left.J_{1}^{\perp}\right)^{\prime}(0)$ is the perpendicular part of the time derivative of the speed at $y_{1}$ of the geodesic in time 1 from $y_{1}$ to $x$. So from equation (B.7) we deduce the rotation

$$
\begin{equation*}
\left(J_{1}^{\perp}\right)^{\prime}(0) d t=\rho_{S}\left(y_{1}\right) \nabla_{t} N^{D_{t}}\left(y_{1}\right)=-\rho_{S}\left(y_{1}\right) \nabla H^{D_{t}}\left(y_{1}\right) d t \tag{B.24}
\end{equation*}
$$

(in the r.h.s. the gradient corresponds to the tangential gradient on $\partial D_{t}$, recall that $H^{D_{t}}$ is only defined on this hypersurface).

We conclude that the horizontal displacement of $x$ is $J_{1}^{T}(1) d t$

$$
J_{1}^{T}(1) d t=J_{1}^{\perp}(1) d t+\frac{1}{4 \sin ^{2} \theta^{S_{t}}(x)}\left(2\left\langle J_{1}^{\perp}(1), N_{2}^{D_{t}}(x)\right\rangle N_{1}^{D_{t}}(x)\right.
$$

$$
\begin{equation*}
\left.+\left(H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)\right)\left(N_{1}^{D_{t}}(x)+N_{2}^{D_{t}}(x)\right)\right) d t \tag{B.25}
\end{equation*}
$$

where $J_{1}^{\perp}(1)=J\left(1,0,-\rho_{S}\left(y_{1}\right) \nabla H^{D_{t}}\left(y_{1}\right)\right)$. Again the processus $z_{t}$ does not play a role.
To summarize, we have the following result for the evolution of $S_{t}$ :

Theorem B.1. When $D_{t}$ evolves as (B.7)

$$
\begin{equation*}
d \partial D_{t}(y)=N^{D_{t}}(y)\left(H^{D_{t}}(y) d t+d z_{t}\right) \tag{B.26}
\end{equation*}
$$

the regular skeleton $S_{t}$ has the normal evolution (B.23)

$$
\begin{equation*}
d S_{t}^{\perp}(x)=\frac{H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)}{4 \sin ^{2} \theta^{S_{t}}(x)}\left(N_{1}^{D_{t}}(x)-N_{2}^{D_{t}}(x)\right) d t \tag{B.27}
\end{equation*}
$$

and the tangential evolution (B.25) which can be rewritten as

$$
\begin{align*}
& d S_{t}^{T}(x) \\
& =p_{S}\left(J_{1}^{\perp}(1)\right) d t  \tag{B.28}\\
& +\left(-\frac{\left\langle J_{1}^{\perp}(1), N_{1}^{S}(x)\right\rangle}{2 \sin \theta^{S_{t}}(x)}+\frac{H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)}{4 \sin ^{2} \theta^{S_{t}}(x)}\right)\left(N_{1}^{D_{t}}(x)+N_{2}^{D_{t}}(x)\right) d t
\end{align*}
$$

where $p_{S}$ denotes the orthogonal projection on $T S, J_{1}^{\perp}(1)=J\left(1,0,-\rho_{S}\left(y_{1}\right) \nabla H^{D_{t}}\left(y_{1}\right)\right)$, and $y_{1}, y_{2}$ are defined in Figure 2.

Remark B.2. The points $y_{1}$ and $y_{2}$ do not play the same role in Theorem B.1. As formula (B.27) is symmetric in $y_{1}$ and $y_{2}$, formula (B.28) is not. The reason is that if we assume the motion of $y_{1}$ to be normal to the boundary $\partial D_{t}$ and to have speed given by (B.26), the motion of $y_{2}$ has no reason to be normal to the boundary: $J_{2}^{\perp}(0)$ does not vanish.

## Appendix C. Doss-Sussman representation of Itô's EQUATION (2.11)

In this section we adapt the results of [6] to our notations. Let the stochastic mean curvature flow be a solution of :
(C.1) $\forall t \in[0, \tau), \forall x \in C_{t}, \quad d \partial D_{t}(y)=\left(d W_{t}+\frac{1}{2} h^{D_{t}}(y) d t\right) N^{D_{t}}(y)$
where $C_{t}:=\partial D_{t}$, starting at $D_{0}$.
Let $\partial G_{t}$ be a solution of
(C.2) $\quad\left\{\begin{array}{rl}G_{0} & =D_{0} \\ \forall t \in[0, \tilde{\epsilon}), \forall x \in \partial G_{t}, & \partial_{t} x\end{array}=\frac{1}{2} \alpha_{\partial G_{t}, W_{t}}(x) N^{G_{t}}(x)\right.$
for some $\tilde{\epsilon}>0$ small enough, where $\alpha$ is defined by
(C.3) $\quad \forall r>0, \forall D \in \mathcal{D}_{r}, \forall x \in C, \quad \alpha_{C, r}(x):=\frac{1}{2} h^{\Psi(C, r)}\left(\psi_{C, r}(x)\right)$
and $\Psi(C, r)$ is the normal flow starting at $C$ at time $r$.
Similarly to the proof of Theorem 9 from [6], we show that $D_{t}=\Psi\left(G_{t},-W_{t}\right)$ is a solution of the stopped martingale problem associated to the generator $(\mathcal{D}, \widetilde{\mathcal{L}})$ where for $f \in C^{\infty}(M)$ and $\mathbb{F}_{f}(D)=\int_{D} f d \mu$,

$$
\widetilde{\mathcal{L}} \mathbb{F}_{f}(D):=\int_{\partial D}\langle\nabla f, \nu\rangle d \underline{\mu}=\mathbb{F}_{\Delta f}(D)
$$

Recall that the equation (C.2), is in fact a quasiparabolic equation with coefficients that depend on trajectory of the Brownian motion (the meaning is trajectory by trajectory). Similarly to Section 4.1 from [6], we show that the solution of (C.2) have a regularity $C^{1+\frac{\alpha}{2}, 2+\alpha}$, for all $\alpha<1$.

Proposition C.1. Let $\partial G_{t}$ be a solution of (C.2). Then $\partial D_{t}=\Psi\left(\partial G_{t},-W_{t}\right)$ is a solution of (C.1) in the Itô sense.

Proof. Let $x \in \Psi\left(\partial G_{t},-W_{t}\right)$, we have :
(C.4)

$$
\begin{aligned}
& d \Psi\left(\partial G_{t},-W_{t}\right)(x)= \\
& =T_{1} \Psi_{\left(\partial G_{t},-W_{t}\right)}\left(\frac{d}{d t} \partial G_{t}\right)\left(\Psi^{-1}\left(\partial G_{t},-W_{t}\right)(x) d t\right. \\
& -N^{\Psi\left(\partial G_{t}, W_{t}\right)}(x) d W_{t} \\
& =\left(-d W_{t}+\frac{1}{2} h^{\Psi\left(\partial G_{t},-W_{t}\right)}(x) d t\right) N^{\Psi\left(\partial G_{t},-W_{t}\right)}(x),
\end{aligned}
$$

where in the first equality we use the Ito formula, and the fact that $t \mapsto \partial G_{t}$ is $C^{1+\frac{\alpha}{2}}$, $\frac{d^{2}}{d^{2} r} \Psi(x, r)=0$, and in the second equality we used Lemma 13 in [6], i.e. $\partial D_{t}$ is a solution in the Itô form :

$$
\left\{\begin{align*}
d \partial D_{t}(x) & =\left(-d W_{t}+\frac{1}{2} h^{\partial D_{t}}(x) d t\right) \nu_{\partial D_{t}}(x)  \tag{C.5}\\
x & \in \partial D_{t} .
\end{align*}\right.
$$

Proposition C.2. Conversely, if $\partial D_{t}$ is a solution of (C.5) then $\partial G_{t}=\Psi\left(\partial D_{t},-W_{t}\right)$ is a solution of (C.2).

Proof. Let $x \in \partial \Psi\left(\partial D_{t},-W_{t}\right)$

$$
\begin{aligned}
& d \Psi\left(\partial D_{t},-W_{t}\right)(x) \\
& =T_{1} \Psi_{\left(\partial D_{t},-W_{t}\right)}\left(\circ d \partial D_{t}\right)(x)+N^{\Psi\left(\partial D_{t}, W_{t}\right)}(x) d B_{t} \\
& =-T_{1} \Psi_{\left(\partial D_{t},-W_{t}\right)}\left(\left(-d W_{t}+\frac{1}{2} h^{\partial D_{t}} d t\right) N^{\partial D_{t}}\right)(x) \\
& -N^{\Psi\left(\partial D_{t}, W_{t}\right)}(x) d W_{t} \\
& =\left(\frac{1}{2} h^{\partial D_{t}}\left(\Psi^{-1}\left(\partial D_{t}, W_{t}\right)(x)\right) N^{\partial G_{t}}(x) d t\right) \\
& =\frac{1}{2} N^{\Psi\left(\partial G_{t},-W_{t}\right)}\left(\Psi\left(\partial G_{t},-W_{t}\right)(x)\right) N^{\partial G_{t}}(x) d t
\end{aligned}
$$

where we use that the Stratonovich differential is equal to the Itô's one, i.e. $\circ d \partial D_{t}(x)=$ $d \partial D_{t}$. So $\partial G_{t}$ is a solution of (C.2).

By the uniqueness of the solution of (C.2) and the fact that it is adapted to the filtration of $B$ we deduce that the solution of (C.5) is unique and is a strong solution. Similarly we have the uniqueness of the solution of

$$
d \partial D_{t}(x)=\left(d W_{t}+\frac{1}{2} h^{\partial D_{t}}(x) d t-\frac{\mu\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)} d t\right) N^{\partial D_{t}}(x) .
$$

Moreover, since we could also make a change of time in the Itô equation, Equation (2.11) has a unique strong solution.

## Appendix D. Weak semi-group theory in the martingale problem sense

This theory has been developed in several books, see for instance Stroock and Varadhan [19] or Ethier and Kurtz [8]. Here we present a minimal version suitable for our purposes.

Let $V$ be a measurable state space and consider $\Omega$ a set of trajectories from $\mathbb{R}_{+}$to $V$. The canonical coordinates on $\Omega$ are denoted by the $X_{t}$, for $t \geqslant 0$ : for $\omega \in \Omega, X_{t}(\omega)$ is the
position at time $t$ of $\omega$. The set $\Omega$ is endowed with the sigma-field generated by the $X_{t}$, for $t \geqslant 0$. Our first assumption is that the mapping

$$
\Omega \times \mathbb{R}_{+} \ni(\omega, t) \quad \mapsto \quad X_{t}(\omega) \in V
$$

is measurable, which usually means that " $\Omega$ is not too big".
For $t \geqslant 0$, we define

$$
\mathcal{F}_{t}:=\sigma\left(X_{s}: s \in[0, t]\right)
$$

For $t \geqslant 0$, we will also need the time shift $\Theta_{t}$ associating to any $\omega \in \Omega$ the trajectory $\Theta_{t}(\omega)$ defined by

$$
\forall s \geqslant 0, \quad X_{s}\left(\Theta_{t}(\omega)\right)=X_{s+t}(\omega)
$$

We assume that $\Theta_{t}(\Omega) \subset \Omega$.
A given family $\mathbb{P}:=\left(\mathbb{P}_{x}\right)_{x \in V}$ of probability measures on $\Omega$ is said to be Markovian if for any $x \in V$ and any $t \geqslant 0$, the image by $\Theta_{t}$ of $\mathbb{P}_{x}$ conditioned by $\mathcal{F}_{t}$ is $\mathbb{P}_{X_{t}}$. In particular, it is assumed that $\mathbb{P}$ has the regularity of a Markov kernel from $V$ to $\Omega$.

From now on, we suppose that a Markovian family $\mathbb{P}$ is given. Let $\mathcal{B}$ be the space of bounded and measurable functions defined on $V$. The semi-group $P:=\left(P_{t}\right)_{t \geqslant 0}$ associated to $\mathbb{P}$ is the family of operators acting on $\mathcal{B}$ via

$$
\forall t \geqslant 0, \forall f \in \mathcal{B}, \forall x \in V, \quad P_{t}[f](x) \quad:=\mathbb{E}_{x}\left[f\left(X_{t}\right)\right]
$$

The Markovianity of $\mathbb{P}$ implies at once the semi-group property

$$
\forall s, t \geqslant 0, \quad P_{t} P_{s}=P_{t+s}
$$

and in particular the elements of $P$ commute.
A subclass of "regular" functions that will be important for our purposes is $\mathcal{R}$ defined as

$$
\mathcal{R}:=\left\{f \in \mathcal{B}: \forall x \in V, \lim _{t \rightarrow 0_{+}} P_{t}[f](x)=f(x)\right\}
$$

Exceptionally in the above limit, we assumed that $t \geqslant 0$ (i.e. not only that $t>0$ ), so that by definition, for any $f \in \mathcal{R}$ and $x \in V, P_{0}[f](x)=f(x)$.

Let us observe that $\mathcal{R}$ is left stable by the semi-group:
Lemma D.1. For any $t \geqslant 0$, we have $P_{t}[\mathcal{R}] \subset \mathcal{R}$. Thus for any given $f \in \mathcal{R}$ and $x \in V$, the mapping

$$
\mathbb{R}_{+} \ni t \quad \mapsto \quad P_{t}[f](x)
$$

is right continuous.
Proof. Indeed, fix $t \geqslant 0$ and $f \in \mathcal{R}$, we have for any $x \in V$ and $s \geqslant 0$,

$$
\begin{aligned}
P_{s}\left[P_{t}[f]\right](x) & =P_{t}\left[P_{s}[f]\right](x) \\
& \left.=\mathbb{E}_{x}\left[P_{s}[f]\left(X_{t}\right)\right]\right]
\end{aligned}
$$

We have for any $s \geqslant 0,\left\|P_{s}[f]\right\|_{\infty} \leqslant\|f\|_{\infty}$ (where $\|\cdot\|_{\infty}$ stands for the supremum norm on $\mathcal{B}$ ) and since $f \in \mathcal{R}$, we get everywhere

$$
\lim _{s \rightarrow 0_{+}} P_{s}[f]\left(X_{t}\right)=f\left(X_{t}\right)
$$

Dominated convergence implies that

$$
\begin{aligned}
\left.\lim _{s \rightarrow 0_{+}} \mathbb{E}_{x}\left[P_{s}[f]\left(X_{t}\right)\right]\right] & =\mathbb{E}_{x}\left[f\left(X_{t}\right)\right] \\
& =P_{t}[f]
\end{aligned}
$$

as desired.
The generator $L$ associated to $P$ is the operator

$$
L: \mathcal{D}(L) \quad \rightarrow \quad \mathcal{R}
$$

defined in the following way: the space $\mathcal{D}(L)$ is the set of functions $f \in \mathcal{R}$ for which there exists a function $g \in \mathcal{R}$ such that the process $M^{f, g}:=\left(M_{t}^{f, g}\right)_{t \geqslant 0}$ defined by

$$
\forall t \geqslant 0, \quad M_{t}^{f, g} \quad:=\quad f\left(X_{t}\right)-f\left(X_{0}\right)-\int_{0}^{t} g\left(X_{s}\right) d s
$$

is a martingale under $\mathbb{P}_{x}$, for all $x \in V$.
Let us remark that $g$ is then uniquely determined. Indeed, we have for any $x \in V$ and $t \geqslant 0$,

$$
\mathbb{E}_{x}\left[f\left(X_{t}\right)\right]-\mathbb{E}\left[f\left(X_{0}\right)\right]-\mathbb{E}\left[\int_{0}^{t} g\left(X_{s}\right) d s\right]=0
$$

Using Fubini's lemma (applicable due to our measurability requirement on $\Omega$ ) and taking into account the definition of $P$, we get

$$
P_{t}[f](x)-P_{0}[f](x)-\int_{0}^{t} P_{s}[g](x) d s=0
$$

namely, recalling that we required that $g \in \mathcal{R}$,

$$
\begin{align*}
g & =P_{0}[g] \\
& =\lim _{t \rightarrow 0_{+}} \frac{1}{t} \int_{0}^{t} P_{s}[g](x) d s \\
& =\lim _{t \rightarrow 0_{+}} \frac{P_{t}[f](x)-f(x)}{t} \tag{D.1}
\end{align*}
$$

(we came back to the usual convention that $t>0$ in the above limit) and as a by-product, we are assured of the existence of the latter limit.

We define $L[f]:=g$ and $M^{f}:=M^{f, g}$.
The differentiation property (D.1) can be extended into
Lemma D.2. For any $f \in \mathcal{D}(L), x \in V$ and $t \geqslant 0$, we have

$$
\begin{equation*}
\partial_{t} P_{t}[f](x)=P_{t}[L[f]](x) \tag{D.2}
\end{equation*}
$$

Proof. For any $f \in \mathcal{D}(L), x \in V$ and $t, s \geqslant 0$, we have

$$
\begin{aligned}
\mathbb{E}_{x}\left[M_{t+s}^{f}-M_{t}^{f}\right] & =\mathbb{E}_{x}\left[\mathbb{E}_{x}\left[M_{t+s}^{f}-M_{t}^{f} \mid \mathcal{F}_{t}\right]\right] \\
& =0
\end{aligned}
$$

We compute that

$$
M_{t+s}^{f}-M_{t}^{f}=f\left(X_{t+s}\right)-f\left(X_{t}\right)-\int_{t}^{t+s} L[f]\left(X_{u}\right) d u
$$

so that

$$
\mathbb{E}_{x}\left[M_{t+s}^{f}-M_{t}^{f}\right]=P_{t+s}[f](x)-P_{t}[f](x)-\int_{0}^{s} P_{t+u}[L[f]](x) d u
$$

Since $L[f] \in \mathcal{R}$, the mapping $[0, s] \ni u \mapsto P_{t+u}[L[f]](x)$ is right continuous, according to Lemma D.1, and the same argument as in (D.1) enables to conclude to (D.2).

We can now come to the main goal of this appendix:
Proposition D.3. For any $t \geqslant 0, \mathcal{D}(L)$ is stable by $P_{t}$ and on $\mathcal{D}(L)$ we have $L P_{t}=P_{t} L$.
Proof. Fix $f \in \mathcal{D}(L)$ and $x \in V$, the assertion of the lemma amounts to checking that the process $N:=\left(N_{s}\right)_{s \geqslant 0}$ defined by

$$
\left(N_{s}\right)_{s \geqslant 0}:=\left(P_{t}[f]\left(X_{s}\right)-P_{t}[f]\left(X_{0}\right)-\int_{0}^{s} P_{t}[L[f]]\left(X_{u}\right) d u\right)_{s \geqslant 0}
$$

is a martingale under $\mathbb{P}_{x}$. Consider $s^{\prime} \geqslant s \geqslant 0$, we have to prove that

$$
\begin{equation*}
\mathbb{E}_{x}\left[N_{s^{\prime}}-N_{s} \mid \mathcal{F}_{s}\right]=0 \tag{D.3}
\end{equation*}
$$

The l.h.s. is equal to

$$
\begin{aligned}
& \mathbb{E}_{x}\left[P_{t}[f]\left(X_{s^{\prime}}\right)-P_{t}[f]\left(X_{s}\right)-\int_{s}^{s^{\prime}} P_{t}[L[f]]\left(X_{u}\right) d u \mid \mathcal{F}_{s}\right] \\
& \quad=\mathbb{E}_{x}\left[P_{t}[f]\left(X_{s^{\prime}-s} \circ \Theta_{s}\right)-P_{t}[f]\left(X_{0} \circ \Theta_{s}\right)-\int_{0}^{s^{\prime}-s} P_{t}[L[f]]\left(X_{u} \circ \Theta_{s}\right) d u \mid \mathcal{F}_{s}\right] \\
& \quad=\mathbb{E}_{y}\left[P_{t}[f]\left(X_{s^{\prime}-s}\right)-P_{t}[f]\left(X_{0}\right)-\int_{0}^{s^{\prime}-s} P_{t}[L[f]]\left(X_{u}\right) d u\right]
\end{aligned}
$$

where $y=X_{s}$. By Fubini's lemma, the previous r.h.s. can be written

$$
\begin{aligned}
\mathbb{E}_{y} & {\left[P_{t}[f]\left(X_{s^{\prime}-s}\right)\right]-\mathbb{E}_{y}\left[P_{t}[f]\left(X_{0}\right)\right]-\int_{0}^{s^{\prime}-s} \mathbb{E}_{y}\left[P_{t}[L[f]]\left(X_{u}\right)\right] d u } \\
& =P_{t+s^{\prime}-s}[f](y)-P_{t}[f](y)-\int_{0}^{s^{\prime}-s} P_{t+u}[L[f]](y) d u
\end{aligned}
$$

Taking into account (D.2), the last integral is equal to

$$
\int_{0}^{s^{\prime}-s} \partial_{u} P_{t+u}[f](y) d u=P_{t+s^{\prime}-s}[f](y)-P_{t}[f](y)
$$

which ends the proof of (D.3).
The advantage of the above approach is that it is quite sable by optional stopping, as it is the case for martingales. Let us succinctly give a simple example in the spirit of Section 2.

Assume that in the above framework, $V$ is a metric space, endowed with its Borelian measurable structure, and that $\Omega$ is the set of continuous trajectories $\mathcal{C}\left(\mathbb{R}_{+}, V\right)$. Furthermore, we suppose that $P$ is Fellerian, in the sense that it preserves $\mathcal{C}_{\mathrm{b}}(V)$, the set of bounded and continuous real functions on $V$.

Let be given $A \subset V$ a closed set. We consider $\tau$ the hitting time of $A$ :

$$
\tau:=\inf \left\{t \geqslant 0: X_{t} \in A\right\} \in \mathbb{R}_{+} \sqcup\{+\infty\}
$$

Define the "new" process $\tilde{X}:=\left(\tilde{X}_{t}\right)_{t \geqslant 0}$ via

$$
\forall t \geqslant 0, \quad \widetilde{X}_{t}:=X_{t \wedge \tau}
$$

and for $x \in V$, let $\widetilde{\mathbb{P}}_{x}$ be the image of $\mathbb{P}_{x}$ by $\widetilde{X}$, it is still a probability measure on $\mathcal{C}\left(\mathbb{R}_{+}, V\right)$. All notions corresponding to $\widetilde{\mathbb{P}}:=\left(\widetilde{\mathbb{P}}_{x}\right)_{x \in V}$, which is still a Markovian family, receive a tilde. It appears without difficulty that $\widetilde{\mathcal{R}}$ is the set of functions $\widetilde{f} \in \mathcal{B}$ such that there exists
$f \in \mathcal{R}$ with $\tilde{f}$ coinciding with $f$ on $V \backslash A$. The domain $\mathcal{D}(\tilde{L})$ is the set of $\tilde{f} \in \widetilde{\mathcal{R}}$ such that there exists $f \in \mathcal{D}(L)$ with $\tilde{f}$ coinciding with $f$ on $V \backslash A$. In addition, we have

$$
\forall x \in V, \quad \widetilde{L}[\tilde{f}](x)= \begin{cases}L[f](x) & , \text { when } x \notin A \\ 0 & , \text { when } x \in A\end{cases}
$$

This expression does not depend on the choice of $f$, due to the fact that $\mathbb{P}$ is a diffusion, i.e. that $\Omega=\mathcal{C}\left(\mathbb{R}_{+}, V\right)$, which implies that $L$ is a local operator (see for instance Theorem 7.29 of Schilling and Partzsch [18], they are working with Euclidean spaces, but the result can be extended to metric spaces).

According to (D.2) and Proposition D.3, we get

$$
\forall \tilde{f} \in \mathcal{D}(\widetilde{L}), \forall x \in V, \forall t \geqslant 0 \quad \partial_{t} \widetilde{P}_{t}[\tilde{f}](x)=\widetilde{P}_{t}[\tilde{L}[\tilde{f}]](x)=\widetilde{L}\left[\widetilde{P}_{t}[\tilde{f}]\right](x)
$$

Such relations are not so obvious if we had chosen to work in a Banach setting (cf. e.g. the book of Yosida [20]), considering for instance semi-groups acting on the space $\mathcal{C}_{\mathrm{b}}(V)$ (endowed with the supremum norm), since in general $\widetilde{L}$ would not naturally take values in $\mathcal{C}_{\mathrm{b}}(V)$.

## Appendix E. A measure theory result

This appendix is not used in this paper. However it could offer an alternative to Lemma H. 1 if we were able to establish that for $i=1,2$ and

$$
\begin{equation*}
\left(\left(X_{t}^{i, \delta}, D_{t}^{\delta}, W_{t}^{i, \delta}, \widetilde{W}_{t}^{\delta}, W_{t}^{i, \delta, m}\right)_{t \geqslant 0}, \tau_{\varepsilon}^{\delta}\right) \tag{E.1}
\end{equation*}
$$

as in (3.18) such that

$$
\left(X_{t}^{i, \delta}, W_{t}^{i, \delta}, W_{t}^{i, \delta, m}\right)_{t \geqslant 0}, \quad i=1,2
$$

are conditionally independent given $\left(\left(D_{t}^{\delta}, \widetilde{W}_{t}^{\delta}\right)_{t \geqslant 0}, \tau_{\varepsilon}^{\delta}\right)$, then the conditional independence remains true in the limit. Conditioning with respect to $\left(\left(D_{t}^{\delta}, \widetilde{W}_{t}^{\delta}\right)_{t \geqslant 0}, \tau_{\varepsilon}^{\delta}\right)$ would allow to work with finite dimensional processes and Lemma 4 in [21] would be sufficient for the convergences required by Theorem 3.5, instead of resorting to Lemma H.1.

Consider $V$ and $W$ two Polish spaces. Let $\mu$ be a probability measure on $V$ and $\left(K_{n}\right)_{n \in \mathbb{N} \sqcup\{\infty\}}$ a family of Markov kernels from $V$ to $W$. For any $n \in \mathbb{N} \sqcup\{\infty\}$, define the probability measure $m_{n}$ on $V \times W \times W$ via

$$
\forall(x, y, z) \in V \times W \times W, \quad m_{n}(d x, d y, d z) \quad:=\mu(d x) K_{n}(x, d y) K_{n}(x, d z)
$$

We have the following result:
Proposition E.1. Assume that the sequence $\left(m_{n}\right)_{n \in \mathbb{N}}$ weakly converges toward $m_{\infty}$ on $V \times W \times W$. Then there exists a subsequence $\left(n_{l}\right)_{l \in \mathbb{N}}$ such that $\mu$-a.s. in $x \in V$, the sequence $K_{n_{l}}(x, \cdot)$ weakly converges toward $K_{\infty}(x, \cdot)$ on $W$.
Proof. Fix $f \in \mathcal{C}_{\mathrm{b}}(V)$ and $g \in \mathcal{C}_{\mathrm{b}}(W)$, where $\mathcal{C}_{\mathrm{b}}(V)$ (respectively $\mathcal{C}_{\mathrm{b}}(W)$ ) stands for the space of bounded continuous functions on $V$ (resp. $W$ ).

Consider the function $h:=f \otimes g \otimes \mathbb{1}_{W}$, where $\mathbb{1}_{W}$ is the mapping on $W$ only taking the value 1 . We have $h \in \mathcal{C}_{\mathrm{b}}\left(V \times W^{2}\right)$, so we have

$$
\lim _{n \rightarrow \infty} m_{n}[h]=m_{\infty}[h]
$$

namely

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \int f K_{n}[g] d \mu=\int f K_{\infty}[g] d \mu \tag{E.2}
\end{equation*}
$$

Let us extend this convergence to any $f \in \mathbb{L}^{2}(\mu)$.
Indeed, given $\epsilon>0$, we can find $\widetilde{f} \in \mathcal{C}_{\mathrm{b}}(V)$ such that

$$
\|f-\widetilde{f}\|_{\mathbb{L}^{2}(\mu)} \leqslant \epsilon
$$

For any $n \in \mathbb{N} \sqcup\{\infty\}$, we have

$$
\begin{aligned}
\left|\int f K_{n}[g] d \mu-\int \tilde{f} K_{n}[g] d \mu\right| & \leqslant \int|f-\tilde{f}|\left|K_{n}[g]\right| d \mu \\
& \leqslant\left\|K_{n}[g]\right\|_{\infty} \int|f-\tilde{f}| d \mu \\
& \leqslant\|g\|_{\infty}\|f-\tilde{f}\|_{\mathbb{L}^{2}(\mu)} \\
& \leqslant\|g\|_{\infty} \epsilon
\end{aligned}
$$

where $\|\cdot\|_{\infty}$ stands for the supremum norm and where we used the Cauchy-Schwarz inequality.

It follows that

$$
\begin{aligned}
\limsup _{n \rightarrow \infty} \int f K_{n}[g] d \mu & \leqslant \epsilon+\lim _{n \rightarrow \infty} \int \tilde{f} K_{n}[g] d \mu \\
& \leqslant \epsilon+\int \tilde{f} K_{\infty}[g] d \mu \\
& \leqslant 2 \epsilon+\int f K_{\infty}[g] d \mu
\end{aligned}
$$

Similarly, we deduce that

$$
\limsup _{n \rightarrow \infty} \int f K_{n}[g] d \mu \geqslant \int f K_{\infty}[g] d \mu-2 \epsilon
$$

Since $\epsilon>0$ can be chosen arbitrary small, we get the validity of (E.2), for any $f \in \mathbb{L}^{2}$, when $g \in \mathcal{C}_{\mathrm{b}}(W)$ is fixed, namely the weak convergence of $\left(K_{n}[g]\right)_{n \in \mathbb{N}}$ toward $K_{\infty}[g]$ in $\mathbb{L}^{2}(\mu)$.

To transform this weak convergence into a strong convergence, it is sufficient to check the convergence of the corresponding $\mathbb{L}^{2}(\mu)$ norms, i.e. that

$$
\lim _{n \rightarrow \infty} \int\left(K_{n}[g]\right)^{2} d \mu=\int\left(K_{\infty}[g]\right)^{2} d \mu
$$

This is also a consequence of the weak convergence of $\left(m_{n}\right)_{n \in \mathbb{N}}$ toward $m$, by considering the mapping $\mathbb{1}_{V} \otimes g \otimes g$ on $V \times W \times W$.

Thus we have shown that for any given $g \in \mathcal{C}_{\mathrm{b}}(W)$, the sequence $\left(K_{n}[g]\right)_{n \in \mathbb{N}}$ converges toward $K_{\infty}[g]$ in $\mathbb{L}^{2}(\mu)$. As a consequence, there exists a subsequence $\left(n_{l}\right)_{l \in \mathbb{N}}$ such that $\left(K_{n_{l}}[g]\right)_{l \in \mathbb{N}}$ converges $\mu$-a.s. toward $K_{\infty}[g]$.

A priori, this subsequence $\left(n_{l}\right)_{l \in \mathbb{N}}$ may depend on $g$, so we resort to a diagonal procedure to avoid this difficulty. More precisely, let $\left(g_{k}\right)_{k \in \mathbb{N}}$ be a sequence of functions from $\mathcal{C}_{\mathrm{b}}(W)$ characterizing the weak convergence on $W$. According to the above arguments, for any $k \in \mathbb{N}$, there exists a subsequence $\left(n_{l}^{(k)}\right)_{l \in \mathbb{N}}$ such that $\left(K_{n_{l}^{(k)}}\left[g_{k}\right]\right)_{l \in \mathbb{N}}$ converges $\mu$-a.s. toward $K_{\infty}\left[g_{k}\right]$. Furthermore, these subsequences can be constructed iteratively: first we find $\left(n_{l}^{(1)}\right)_{l \in \mathbb{N}}$, next $\left(n_{l}^{(2)}\right)_{l \in \mathbb{N}}$ is obtained as a subsequence of $\left(n_{l}^{(1)}\right)_{l \in \mathbb{N}}$, and so on, for any
$k \in \mathbb{N},\left(n_{l}^{(k+1)}\right)_{l \in \mathbb{N}}$ is a subsequence of $\left(n_{l}^{(k)}\right)_{l \in \mathbb{N}}$. Define the subsequence $\left(n_{l}\right)_{l \in \mathbb{N}}$ via

$$
\forall l \in \mathbb{N}, \quad n_{l} \quad:=n_{l}^{(l)}
$$

The sequence $\left(n_{l}\right)_{l \in \mathbb{N}}$ is a subsequence of all the subsequences $\left(n_{l}^{(k)}\right)_{l \in \mathbb{N}}$, and so for any $k \in \mathbb{N},\left(K_{n_{l}}\left[g_{k}\right]\right)_{l \in \mathbb{N}}$ converges $\mu$-a.s. toward $K_{\infty}\left[g_{k}\right]$. Taking the union of the underlying $\mu$-negligible sets, we get that $\mu$-a.s., for any $k \in \mathbb{N},\left(K_{n_{l}}\left[g_{k}\right]\right)_{l \in \mathbb{N}}$ converges toward $K_{\infty}\left[g_{k}\right]$. By choice of the sequence $\left(g_{k}\right)_{k \in \mathbb{N}}$, the desired result follows.

## Appendix F. An Itô-TANAKA FORMULA

Let $M$ be a $d$-dimensional Riemannian manifold and $D \subset M$ a compact and connected domain with $C^{2}$ boundary $\partial D$, and $S$ be the regular skeleton of $D$, and $\rho_{\partial D}^{+}$the signed distance to $\partial D$, which is positive inside $D$ and negative outside $D$. The notations will be the same as in Appendix A.

Proposition F.1. Let $X_{t}$ a Brownian motion in M. We have the following Itô-Tanaka formula:

$$
d \rho_{\partial D}^{+}\left(X_{t}\right)=\left\langle N^{D}\left(X_{t}\right), d X_{t}\right\rangle-\frac{1}{2} h^{D}\left(X_{t}\right) d t-\sin \left(\theta^{S}\left(X_{t}\right)\right) d L_{t}^{S}(X)
$$

in the above formula, $N^{D}(x)=\nabla \rho_{\partial D}^{+}(x)$ and $-h^{D}(x)=\Delta \rho_{\partial D}^{+}(x)$ for $x \notin S$, and define to be 0 elsewhere, $L_{t}^{S}(X)$ is the local time defined as in (3.11).

Proof. The formula is a consequence of the Itô formula outside the skeleton. Since the non regular part of the skeleton has Hausdorff dimension smaller than or equal to $d-2$, it is not visited by the Brownian motion. So we only focus on the regular skeleton. For all $x \in S$, the distance to the boundary is the minimum of two $C^{2}$ functions $f, g$ defined on some neighborhood $U$ of $x$ in $M$. The function $f$ (resp. $g$ ) is the distance function to a piece of $\partial D$ containing $y_{1}$ (resp. $y_{2}$ ) as in (A.2). We have locally,

$$
\rho_{\partial D}^{+}=f \wedge g=\frac{1}{2}(f+g)-\frac{1}{2}|f-g|
$$

Using Itô formula and Tanaka formula we have

$$
\begin{aligned}
d \rho_{\partial D}^{+}\left(X_{t}\right) & =\frac{1}{2}\left(\frac{1}{2} \Delta(f+g)\left(X_{t}\right) d t+\left\langle\nabla(f+g)\left(X_{t}\right), d X_{t}\right\rangle\right) \\
& -\frac{1}{2}\left(\operatorname{sign}\left((f-g)\left(X_{t}\right)\right) d\left((f-g)\left(X_{t}\right)\right)+d L_{t}^{0,+}\left((f-g)\left(X_{.}\right)\right)\right)
\end{aligned}
$$

where $L_{t}^{0,+}((f-g)(X))=.\lim _{\varepsilon \rightarrow 0^{+}} \frac{1}{\varepsilon} \int_{0}^{t} \mathbb{1}_{[0, \varepsilon]}\left((f-g)\left(X_{s}\right)\right) d\langle(f-g)(X),(f-g)(X)\rangle_{s}$. Since locally $S=\{f-g=0\}$ and $\mu(S)=0$, we have
$d \rho_{\partial D}^{+}\left(X_{t}\right)=\frac{1}{2} \mathbb{1}_{X_{t} \notin S} \Delta \rho_{\partial D}^{+}\left(X_{t}\right) d t+\mathbb{1}_{X_{t} \notin S}\left\langle\nabla \rho_{\partial D}^{+}\left(X_{t}\right), d X_{t}\right\rangle-\frac{1}{2} d L_{t}^{0,+}\left((f-g)\left(X_{.}\right)\right)$.

After changing the role of $f$ and $g$ we get
(F.1)

$$
d \rho_{\partial D}^{+}\left(X_{t}\right)=\frac{1}{2} \mathbb{1}_{X_{t} \notin S} \Delta \rho_{\partial D}^{+}\left(X_{t}\right) d t+\mathbb{1}_{X_{t} \notin S}\left\langle\nabla \rho_{\partial D}^{+}\left(X_{t}\right), d X_{t}\right\rangle-\frac{1}{2} d L_{t}^{0}\left((f-g)\left(X_{.}\right)\right)
$$

where

$$
L_{t}^{0}\left((f-g)\left(X_{.}\right)\right)=\lim _{\varepsilon \rightarrow 0^{+}} \int_{0}^{t} \frac{1}{2 \varepsilon} \mathbb{1}_{[-\varepsilon, \varepsilon]}\left((f-g)\left(X_{s}\right)\right)\|\nabla(f-g)\|^{2}\left(X_{s}\right) d s
$$

In Appendix A it is shown that for $x \in S,\|\nabla(f-g)(x)\|=2 \sin \left(\theta^{S}(x)\right)$.
Using the flow $\frac{d}{d t} \gamma(t)=-\frac{\nabla(f-g)(\gamma(t))}{\| \nabla(f-g)\left(\gamma(t) \|^{2}\right.}$ that starts at $y \in U$, we get

$$
\{y \in M, \text { s.t. }|f-g|(y) \leqslant \varepsilon\} \subset\left\{y \in M, \text { s.t. }\left|d_{S}(y)\right| \leqslant \frac{\varepsilon}{2 \sin \left(\theta^{S}(\gamma(g(y)))\right)}+o(\varepsilon)\right\}
$$

where $d_{S}$ is the distance to $S$. On the other hand, using the minimal geodesic from $S$ to $y \in U$ we get

$$
\left\{y \in M, \text { s.t. }\left|d_{S}(y)\right| \leqslant \varepsilon\right\} \subset\left\{y \in M, \text { s.t. }|f-g|(y) \leqslant 2 \varepsilon \sin \left(\theta^{S}\left(P^{S}(y)\right)\right)+o(\varepsilon)\right\}
$$

Hence

$$
d L_{t}^{0}\left((f-g)\left(X_{.}\right)\right)=2 \sin \left(\theta^{S}\left(X_{t}\right)\right) L_{t}^{S}\left(X_{.}\right)
$$

Together with (F.1), this yield the Proposition.

## Appendix G. Uniqueness in Law of $\tilde{\mathcal{L}}$ diffusion

Let us consider the following generator $\widehat{\mathscr{L}}$ of a stochastic modified mean curvature flow. The action of this generator and its carré du champs on elementary observables are defined as follows. For any smooth function $k$ on $M$, consider the mapping $F_{k}$ on $\mathcal{D}^{2+\alpha}$ defined by

$$
\forall D \in \mathcal{D}^{2+\alpha}, \quad F_{k}(D):=\int_{D} k d \mu
$$

For any $k, g \in \mathcal{C}^{\infty}(M)$ and any $D \in \mathcal{D}^{2+\alpha}$,

$$
\left\{\begin{align*}
\widehat{\mathscr{L}}\left[F_{k}\right](D) & :=-\frac{1}{2} \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)=F_{\frac{1}{2} \Delta k}(D)  \tag{G.1}\\
\Gamma_{\widehat{\mathscr{L}}}\left[F_{k}, F_{g}\right](D) & :=\int_{\partial D} k d \underline{\mu} \int_{\partial D} g d \underline{\mu} .
\end{align*}\right.
$$

Note that $\widehat{\mathscr{L}}$ has the same carré du champs as the carré du champs associated to $\widetilde{\mathscr{L}}$. From now the generator $\widehat{\mathscr{L}}$ is defined as in (2.14).

Proposition G.1. The martingale problem associated $\widehat{\mathscr{L}}$ is well-posed.
Proof. We have already shown the existence result in [6], so it remains to prove the uniqueness in law. Let us first consider the two-dimensional Euclidean case, namely $M=\mathbb{R}^{2}$. For all $\lambda \in \mathbb{R}$ and for any function $k_{\lambda} \in \operatorname{vect}\left(e^{\lambda x}, e^{\lambda y}\right)$ we have $\frac{1}{2} \Delta k_{\lambda}(x, y)=\frac{\lambda^{2}}{2} k_{\lambda}(x, y)$.

Let $f_{\lambda}((x, y), D):=k_{\lambda}(x, y) F_{k_{\lambda}}(D)$, for $(x, y) \in \mathbb{R}^{2}$ and $D \in \mathcal{D}^{2+\alpha}$. This function satisfies the following property:

$$
\begin{aligned}
\widehat{\mathscr{L}} f_{\lambda}((x, y), D) & =k_{\lambda}(x, y) \widehat{\mathscr{L}} F_{k_{\lambda}}(D) \\
& =k_{\lambda}(x, y) F_{\frac{1}{2} \Delta k_{\lambda}}(D) \\
& =k_{\lambda}(x, y) F_{\frac{\lambda^{2}}{2} k_{\lambda}}(D) \\
& =\frac{\lambda^{2}}{2} k_{\lambda}(x, y) F_{k_{\lambda}}(D) \\
& =\frac{1}{2} \Delta k_{\lambda}(x, y) F_{k_{\lambda}}(D) \\
& =\frac{1}{2} \Delta f_{\lambda}((x, y), D)
\end{aligned}
$$

Let $\left(X_{t}\right)_{t \geqslant 0}$ be a $\mathbb{R}^{2}$-valued Brownian motion that starts at $X_{0}=\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}$ and $\left(\hat{D}_{t}\right)_{t \geqslant 0}$ a $\widehat{\mathscr{L}}$ diffusion that starts at $D_{0}$ independent of $\left(X_{t}\right)_{t \geqslant 0}$. Even if we stop the diffusion, we can assume that its lifetime is infinite and we add indicators as described in Appendix D. For all $0 \leqslant s \leqslant t$, we have

$$
d f_{\lambda}\left(X_{t-s}, \hat{D}_{s}\right) \stackrel{m}{=}-\frac{1}{2} \Delta f_{\lambda}\left(X_{t-s}, \hat{D}_{s}\right) d s+\widehat{\mathscr{L}} f_{\lambda}\left(X_{t-s}, \hat{D}_{s}\right) d s \stackrel{m}{=} 0
$$

Hence for all $\lambda \in \mathbb{R}$ we have

$$
\begin{equation*}
\mathbb{E}\left[f_{\lambda}\left(X_{t}, D_{0}\right)\right]=\mathbb{E}\left[f_{\lambda}\left(X_{0}, \hat{D}_{t}\right)\right] \tag{G.2}
\end{equation*}
$$

Since the left hand side of the above equation does not depend on the $\widehat{\mathscr{L}}$ diffusion, we get that for any $\widehat{\mathscr{L}}$ diffusion $\left(\tilde{D}_{t}\right)_{t \geqslant 0}$ that starts at $D_{0}$ :

$$
\mathbb{E}\left[f_{\lambda}\left(X_{0}, \hat{D}_{t}\right)\right]=\mathbb{E}\left[f_{\lambda}\left(X_{0}, \tilde{D}_{t}\right)\right]
$$

and so

$$
\left.\mathbb{E}\left[F_{k_{\lambda}}\left(D_{t}\right)\right]=\mathbb{E}\left[F_{k_{\lambda}}\left(\tilde{D}_{t}\right)\right)\right]
$$

In order to apply Theorem 4.2 of [8], we have to show that the above equation characterizes the law of the one-dimensional distribution, i.e. we have to show that $\left(F_{k_{\lambda}}\right)$ is separating in the space of probability measures on $\mathcal{D}^{2+\alpha}$. This is equivalent to separate domains. Let $A, B \in \mathcal{D}^{2+\alpha}$ such that $F_{k_{\lambda}}(A)=F_{k_{\lambda}}(B)$ for all $\lambda \in \mathbb{R}$ and $k_{\lambda} \in\left\langle e^{\lambda x}, e^{\lambda y}\right\rangle$, we have for all $\lambda$ :

$$
\int_{A} k_{\lambda}(x, y) d \mu=\int_{B} k_{\lambda}(x, y) d \mu
$$

After successive derivations in $\lambda$ and evaluation at $\lambda=0$, we get for all $n \in \mathbb{N}$

$$
\begin{aligned}
& \int_{A} x^{n} d \mu=\int_{B} x^{n} d \mu \\
& \int_{A} y^{n} d \mu=\int_{B} y^{n} d \mu
\end{aligned}
$$

The above computations could be done also for $\tilde{k}_{\lambda_{1}, \lambda_{2}}=e^{\lambda_{1} x+\lambda_{2} y}$, since $\frac{1}{2} \Delta \tilde{k}_{\lambda_{1}, \lambda_{2}}=$ $\frac{\lambda_{1}^{2}+\lambda_{2}^{2}}{2} \tilde{k}_{\lambda_{1}, \lambda_{2}}$, and after derivations in $\lambda_{1}, \lambda_{2}$ and evaluating at $(0,0)$ we get that for all $n, m \in \mathbb{N}$ :

$$
\int_{A} x^{n} y^{m} d \mu=\int_{B} x^{n} y^{m} d \mu
$$

hence, using the boundary regularity, we get $A=B$.
We could also apply Stone-Weierstrass' theorem to the function algebra generated by the mappings $(x, y) \mapsto e^{\lambda_{1} x}$ and $(x, y) \mapsto e^{\lambda_{2} y}$.

The proof is the same for all Euclidean spaces.
If $M$ is a compact manifold let

$$
f_{\lambda_{i}}(X, D):=k_{\lambda_{i}}(X) F_{k_{\lambda_{i}}}(D)
$$

where $\lambda_{i}$ is an eigenvalue of $\frac{1}{2} \Delta$ and $k_{i}$ is the associated eigenfunction (respectively the Neumann eigenvalue). By the same computation as above (G.2) is also valid for the boundary reflecting Brownian motion), to get the conclusion we have to show that $\left(F_{k_{\lambda_{i}}}\right)_{i}$ separates domains. Since $\left(k_{\lambda_{i}}\right)_{i}$ is an orthonormal basis of $L^{2}(\mu)$ we get that if $A, B \in \mathcal{D}^{2+\alpha}$ be such that for all $i$,

$$
F_{k_{\lambda_{i}}}(A)=F_{k_{\lambda_{i}}}(B)
$$

i.e $\left\langle\mathbb{1}_{A}, k_{\lambda_{i}}\right\rangle_{L^{2}}=\left\langle\mathbb{1}_{B}, k_{\lambda_{i}}\right\rangle_{L^{2}}$, then $\mathbb{1}_{A} \stackrel{L^{2}}{=} \mathbb{1}_{B}$ hence $A=B$.

For the complete manifold $M$, let $\Omega_{k}$ be an exhaustion of $M$ with a regular boundary such that $D_{0} \subset \Omega_{k}$, and stop the $\widehat{\mathscr{L}}$ diffusion when it hit $\Omega_{k}^{c}$ and use the above result for the manifold with boundary $\Omega_{k}$, we get the result by localization.

Proposition G.2. The martingale problem associated to $\mathscr{L}$ is well-posed.
Proof. Let $D_{t}$ be a $\mathscr{L}$ diffusion that starts at $D_{0}$, defined on $\left(\Omega, \mathcal{F}^{D}, \mathbb{Q}\right)$. We first recall that there exist an enlargement of the probability space such that it carries a one dimensional Brownian motion $B$ such that for all $k \in C^{\infty}(M)$

$$
\begin{equation*}
F_{k}\left(D_{t}\right)=F_{k}\left(D_{0}\right)+\int_{0}^{t} \mathscr{L}\left[F_{k}\right]\left(D_{s}\right) d s+\int_{0}^{t} \sqrt{\Gamma_{\mathscr{L}}\left[F_{k}, F_{k}\right]}\left(D_{s}\right) d B_{s} \tag{G.3}
\end{equation*}
$$

where $\sqrt{\Gamma_{\mathscr{L}}\left[F_{k}, F_{k}\right]}(D):=\int_{\partial D} k d \sigma$, this is actually Proposition 53 in [6]. Note that this procedure of enlargement (Theorem 1.7 chapter V in [17]) could be done by gluing the same independent Brownian motion for each $\left(\Omega, \mathcal{F}^{D}, \mathbb{Q}\right)$. We denote by $\left(\tilde{\Omega}, \tilde{\mathcal{F}}^{D}, \tilde{\mathbb{Q}}\right)$ the enlarged probability space. Since $\mathscr{L}$ is an $h$-transform of $\widehat{\mathscr{L}}$ namely

$$
\mathscr{L}\left[F_{k}\right]=\widehat{\mathscr{L}}\left[F_{k}\right]+\frac{\Gamma_{\widehat{\mathscr{L}}}\left(F_{1}, F_{k}\right)}{F_{1}}
$$

equation (G.3) becomes in a differential form

$$
\begin{equation*}
d F_{k}\left(D_{t}\right)-\widehat{\mathscr{L}}\left[F_{k}\right]\left(D_{t}\right) d t=\left(\int_{\partial D} k d \sigma\right)\left(d B_{t}+\frac{\mu^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)} d t\right) \tag{G.4}
\end{equation*}
$$

Let

$$
\begin{gathered}
\left.M_{t}=e^{-\int_{0}^{t}\left\langle\frac{\hat{\mu}^{\partial D_{s}\left(\partial D_{s}\right)}}{\mu\left(D_{s}\right)}\right.} d B_{s}\right\rangle-\frac{1}{2} \int_{0}^{t}\left(\frac{\mu^{\partial D_{s}\left(\partial D_{s}\right)}}{\mu\left(D_{s}\right)}\right)^{2} d s \\
\mathbb{P}_{\mid \mathcal{F}_{t}}=M_{t} \tilde{\mathbb{Q}}_{\mid \mathcal{F}_{t}}
\end{gathered}
$$

Using Girsanov transform, $D_{t}$ is solution of the $\widehat{\mathscr{L}}$ martingale problem on the probability space $\left(\tilde{\Omega}, \tilde{\mathcal{F}}^{D}, \mathbb{P}\right)$. Since $\tilde{\mathbb{Q}}=M^{-1} \mathbb{P}$ we get the uniqueness in law of the $\mathscr{L}$ diffusion by Proposition G.1.

## Appendix H. Convergence in law: A Key lemma

This Appendix is devoted to the adaptation to some domain-valued sequences of processes, of Lemma 4 in [21], which states stability of some time integrals under convergence in law.
Lemma H.1. Let $\tilde{\mathscr{F}}:=\tilde{\mathscr{F}}^{\alpha, \varepsilon}$. We endow the set of continuous paths $\mathscr{C}([0, \infty), M \times \tilde{\mathscr{F}})$ with the two dissimilarity measures $d_{\beta}, \beta \in\{0, \alpha\}$, defined as:

$$
\begin{equation*}
d_{\beta}\left(\left(x^{1}, D^{1}\right),\left(x^{2}, D^{2}\right)\right)=\sup _{t \geqslant 0} \rho\left(x^{1}(t), x^{2}(t)\right)+\sup _{t \geqslant 0} d_{\beta, \tilde{\mathscr{F}}}\left(D^{1}(t), D^{2}(t)\right) \tag{H.1}
\end{equation*}
$$

where for two domains $D$ and $D^{\prime}$

$$
d_{\beta, \tilde{\mathscr{F}}}\left(D, D^{\prime}\right)=\left\{\begin{array}{cc}
d_{\beta, D}\left(D, D^{\prime}\right) \wedge d_{\beta, D^{\prime}}\left(D^{\prime}, D\right) \wedge \varepsilon & \text { if } H\left(D, D^{\prime}\right)<\varepsilon  \tag{H.2}\\
\varepsilon & \text { otherwise. }
\end{array}\right.
$$

Here $H\left(D, D^{\prime}\right)$ is the Hausdorff distance between $D$ and $D^{\prime}$ and the distance $d_{\beta, D}$ is defined in (2.2).

Let $\left(X_{t}^{n}, D_{t}^{n}, \tau_{\varepsilon}^{n}\right)_{t \geqslant 0}:=\left(X_{t}^{\delta_{n}}, D_{t}^{\delta_{n}}, \tau_{\varepsilon}^{\delta_{n}}\right)_{t \geqslant 0}$ a subsequence of (3.18) converging in law to the limit defined in $(3.19)$ for the product of $d_{\alpha}$ and the Euclidean distance in $\mathbb{R}_{+}$.

Let $f_{n}:(x, D) \mapsto f_{n}(x, D)$ and $f:(x, D) \mapsto f(x, D)$ be maps on $M \times \tilde{\mathscr{F}}$ with values in some Euclidean space, and $U$ an open set in $M \times \tilde{\mathscr{F}}$ for $d_{0}$. Assume that:
(i) the random variables $\int_{0}^{\infty}\left|f_{n}\left(X_{s}^{n}, D_{s}^{n}\right)\right|^{p} d s$ are uniformly bounded in probability for some $p>1$,
(ii) in the open set $U$, the functions $f_{n}$ converge locally uniformly to $f$ with respect to $d_{0}$, and are $d_{0}$-continuous,
(iii) for a.e. $t \geqslant 0,\left(X_{t}, D_{t}\right) \in U$.

Then $\left(X_{t}^{n}, D_{t}^{n}, \int_{0}^{t} f_{n}\left(X_{s}^{n}, D_{s}^{n}\right) d s\right)_{t \geqslant 0}$ converges in law to $\left(X_{t}, D_{t}, \int_{0}^{t} f\left(X_{s}, D_{s}\right) d s\right)_{t \geqslant 0}$ for $\left(d_{\alpha},|\cdot|\right)$.
Remark H.2. In the applications we will always take

$$
\begin{equation*}
U=\{(x, D) \in M \times \tilde{\mathscr{F}}, x \in D \backslash S(D)\} \tag{H.3}
\end{equation*}
$$

which is easily seen to be $d_{0}$-open thanks to Assumption 3.1 on $\tilde{\mathscr{F}}$.
Proof. We will follow the proof of Lemma 4 in [21], but with several differences due to infinite dimensional spaces. Set for $n \in \mathbb{N}, t \geqslant 0$,

$$
\begin{equation*}
A_{t}^{n}:=\int_{0}^{t} f_{n}\left(X_{s}^{n}, D_{s}^{n}\right) d s, \quad A_{t}:=\int_{0}^{t} f\left(X_{s}, D_{s}\right) d s \tag{H.4}
\end{equation*}
$$

Condition (i) implies that the processes $A^{n}$ are tight. To get the conclusion il is sufficient to show that all the converging subsequences have the same limit. So assume that

$$
\begin{equation*}
\left(X_{t}^{n}, D_{t}^{n}, A_{t}^{n}\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(X_{t}, D_{t}, a_{t}\right)_{t \geqslant 0} . \tag{H.5}
\end{equation*}
$$

and let us prove that $\left(a_{t}\right)_{t \geqslant 0}=\left(A_{t}\right)_{t \geqslant 0}$. By Skorohod theorem we may realize all processes

$$
\begin{equation*}
\left(X_{t}^{n}, D_{t}^{n}, A_{t}^{n}, X_{t}, D_{t}, a_{t}\right)_{t \geqslant 0} \tag{H.6}
\end{equation*}
$$

on the same probability space $(\Omega, \mathscr{F}, \mathbb{P})$ in such a way that

$$
\begin{equation*}
\left(Z_{t}^{n}\right)_{t \geqslant 0}:=\left(X_{t}^{n}, D_{t}^{n}, A_{t}^{n}\right)_{t \geqslant 0} \xrightarrow{\text { a.s. }}\left(X_{t}, D_{t}, a_{t}\right)_{t \geqslant 0}=:\left(Z_{t}\right)_{t \geqslant 0} \tag{H.7}
\end{equation*}
$$

This means that $Z_{t}^{n} \rightarrow Z_{t}$ a.s. uniformly in $t \geqslant 0$.
Fix $\omega \in \Omega$. Let $t>0$ be such that $\left(X_{t}(\omega), D_{t}(\omega)\right) \in U$. For some $\varepsilon^{\prime}>0$ we have $\left(X_{s}(\omega), D_{s}(\omega)\right) \in U$ for all $s \in\left[t-\varepsilon^{\prime}, t+\varepsilon^{\prime}\right]$. The set

$$
\begin{equation*}
S:=\left\{\left(X_{s}(\omega), D_{s}(\omega)\right), \quad s \in\left[t-\varepsilon^{\prime}, t+\varepsilon^{\prime}\right]\right\} \tag{H.8}
\end{equation*}
$$

is $d_{\alpha}$-compact in $M \times \tilde{\mathscr{F}}$, so it has a $d_{\alpha}$-neighbourhood $V$ included in $U$ of the form

$$
\begin{equation*}
V=\left\{(x, D) \in M \times \tilde{\mathscr{F}}, \quad d_{\alpha}((x, D), S) \leqslant \varepsilon^{\prime \prime}\right\} \tag{H.9}
\end{equation*}
$$

for some small enough $\varepsilon^{\prime \prime}>0$. For $n$ sufficiently large, $\left(X_{s}^{n}(\omega), D_{s}^{n}(\omega)\right) \in V$ for all $s \in$ $\left[t-\varepsilon^{\prime}, t+\varepsilon^{\prime}\right]$. On the other hand $V$ is bounded for the distance $d_{\alpha}$. This implies by ArzelaAscoli theorem that it is compact for the distance $d_{0}$. We have the two following facts, the first one being an assumption on the $f_{n}$ and $f$, the second one being a consequence of the $d_{0}$-compactness of $V$
(a) $f_{n} \rightarrow f$ as $n \rightarrow \infty$ uniformly in $\left(V, d_{0}\right)$;
(b) $f$ is uniformly continuous in $\left(V, d_{0}\right)$.

Then

$$
\begin{aligned}
& \sup _{s \in[t-\varepsilon, t+\varepsilon]}\left|f_{n}\left(X_{s}^{n}(\omega), D_{s}^{n}(\omega)\right)-f\left(X_{s}(\omega), D_{s}(\omega)\right)\right| \\
& \leqslant \sup _{s \in[t-\varepsilon, t+\varepsilon]}\left|f_{n}\left(X_{s}^{n}(\omega), D_{s}^{n}(\omega)\right)-f\left(X_{s}^{n}(\omega), D_{s}^{n}(\omega)\right)\right| \\
& \quad+\sup _{s \in[t-\varepsilon, t+\varepsilon]}\left|f\left(X_{s}^{n}(\omega), D_{s}^{n}(\omega)\right)-f\left(X_{s}(\omega), D_{s}(\omega)\right)\right| .
\end{aligned}
$$

Both terms in the right converge to 0 , the first one by (a) and the second one by (b). So we have by (H.7) and the above calculation
(H.10)

$$
\left\{\begin{aligned}
\left(A_{s}^{n}(\omega)\right)_{s \in[t-\varepsilon, t+\varepsilon]} & \rightarrow\left(a_{s}(\omega)\right)_{s \in[t-\varepsilon, t+\varepsilon]} \\
\left(\left(A_{s}^{n}(\omega)\right)^{\prime}=f_{n}\left(X_{s}^{n}(\omega), D_{s}^{n}(\omega)\right)\right)_{s \in[t-\varepsilon, t+\varepsilon]} & \rightarrow\left(f\left(X_{s}(\omega), D_{s}(\omega)\right)\right)_{s \in[t-\varepsilon, t+\varepsilon]}
\end{aligned}\right.
$$

both uniformly in $s \in[t-\varepsilon, t+\varepsilon]$. This implies that $a_{s}(\omega)$ is differentiable in $(t-\varepsilon, t+\varepsilon)$ with derivative $f\left(X_{s}(\omega), D_{s}(\omega)\right)$ and in particular at $t$.

We have that for all $t \geqslant 0,\left(X_{t}(\omega), D_{t}(\omega)\right) \in U$ a.s.. So for all $t \geqslant 0$,

$$
\begin{equation*}
\frac{d}{d t} a_{t}(\omega)=f\left(X_{t}(\omega), D_{t}(\omega)\right) \quad \text { a.s.. } \tag{H.11}
\end{equation*}
$$

This implies that $\omega$ a.s.

$$
\begin{equation*}
\frac{d}{d t} a_{t}(\omega)=f\left(X_{t}(\omega), D_{t}(\omega)\right) \quad \text { for a.e. } t \tag{H.12}
\end{equation*}
$$

On the other hand we know by [12] Theorem 10 that $\left(a_{t}\right)_{t \geqslant 0}$ is absolutely continuous :

$$
\begin{equation*}
a_{t}(\omega)=\int_{0}^{t} \ell_{s}(\omega) d s \tag{H.13}
\end{equation*}
$$

By Lebesgue theorem, $\omega$ a.s., for a.e. $t \geqslant 0$

$$
\begin{equation*}
\lim _{\varepsilon \searrow 0} \frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon}\left|\ell_{s}(\omega)-\ell_{t}(\omega)\right| d s=0 \tag{H.14}
\end{equation*}
$$

Equalities (H.12) and (H.13) imply that $\omega$ a.s.

$$
\begin{equation*}
\lim _{\varepsilon \searrow 0} \frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon} \ell_{s}(\omega) d s=f\left(X_{t}(\omega), D_{t}(\omega)\right) \quad \text { for a.e. } t \tag{H.15}
\end{equation*}
$$

On the other hand

$$
\left|\frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon} \ell_{s}(\omega)-\ell_{t}(\omega) d s\right| \leqslant \frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon}\left|\ell_{s}(\omega)-\ell_{t}(\omega)\right| d s
$$

so (H.14) implies that $\omega$ a.s. for a.e. $t \geqslant 0$

$$
\begin{equation*}
\lim _{\varepsilon \searrow 0} \frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon} \ell_{s}(\omega) d s=\ell_{t}(\omega) \tag{H.16}
\end{equation*}
$$

Consequently, using (H.12) and (H.16), we get $\omega$ a.s. for a.e. $t \geqslant 0$

$$
\begin{equation*}
\ell_{t}(\omega)=f\left(X_{t}(\omega), D_{t}(\omega)\right) \tag{H.17}
\end{equation*}
$$

Integrating we get $\omega$-a.s. for all $t \geqslant 0$

$$
\begin{equation*}
a_{t}(\omega)=A_{t}(\omega)=\int_{0}^{t} f\left(X_{s}(\omega), D_{s}(\omega)\right) d s \tag{H.18}
\end{equation*}
$$

This together with (H.4) proves the lemma.

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