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MERSENNE

# COUPLINGS OF BROWNIAN MOTIONS WITH SET-VALUED DUAL PROCESSES ON RIEMANNIAN MANIFOLDS 

by Marc Arvaudon, Koléhè Coulibaly-Pasquier<br>\& Laurent Miclo


#### Abstract

The purpose of this paper is to construct a Brownian motion $\left(X_{t}\right)_{t \geqslant 0}$ taking values in a Riemannian manifold $M$, together with a compact set-valued process $\left(D_{t}\right)_{t \geqslant 0}$ such that, at least for small enough $\mathscr{F}^{D}$-stopping time $\tau>0$ and conditioned by $\mathscr{F}_{\tau}^{D}$, the law of $X_{\tau}$ is the normalized Lebesgue measure on $D_{\tau}$. This intertwining result is a generalization of Pitman's 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 each $D_{t}$ play an important role. Several examples with moving intervals, discs, annuli, symmetric convex sets are investigated. Résumé (Couplage des mouvements browniens avec des processus duaux à valeurs ensembles sur des variétés riemanniennes)

L'objectif de cet article est de construire un mouvement brownien $\left(X_{t}\right)_{t \geqslant 0}$ à valeurs dans une variété riemannienne $M$ conjointement avec un processus à valeurs ensembles $\left(D_{t}\right)_{t \geqslant 0}$, de telle sorte qu'au moins pour tout temps d'arrêt $\tau>0$ assez petit dans la filtration $\mathcal{F}^{D}$ engendrée par $\left(D_{t}\right)_{t \geqslant 0}$, la loi de $X_{\tau}$ conditionnée par $\mathcal{F}_{\tau}^{D}$ est la mesure riemannienne conditionnée sur $D_{\tau}$. Ce résultat d'entrelacement est une généralisation du théorème de Pitman. Nous commençons par construire des processus entrelacés réguliers par le biais du théorème de Stokes. Puis en utilisant différentes procédures de limites, nous construisons des processus entrelacés synchrones, libres et miroirs. Les temps locaux du mouvement brownien sur le squelette (morphologique) ou sur la frontière jouent des rôles importants. Nous étudions plusieurs exemples consistant en des intervalles, des disques, des anneaux et des ensembles convexes symétriques.


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## 1. Introduction and main results

Markov intertwinings were introduced by Rogers and Pitman [20] to give a direct proof of the famous relation between the Brownian motion and the Bessel-3 process due to Pitman [18]. These relations were next used by Yor and his coauthors (see e.g. [23, 6]) to get identities in law and by Diaconis and Fill [9] to construct strong stationary times. For a historical account of the subsequent development of the Markov intertwining technique, consult for instance Pal and Shkolnikov [17].

At an algebraic level, a Markov intertwining relation is a (directed) weak similar relation, from a Markov semi-group $\left(\bar{P}_{t}\right)_{t \geqslant 0}$ on a measurable state space $(\bar{M}, \overline{\mathcal{M}})$ to another Markov semi-group $\left(P_{t}\right)_{t \geqslant 0}$ on a measurable state space $(M, \mathcal{M})$, consisting of a Markov kernel (called the link) $\Lambda$ from $(\bar{M}, \overline{\mathcal{M}})$ to $(M, \mathcal{M})$ such that

$$
\begin{equation*}
\forall t \geqslant 0, \quad \bar{P}_{t} \Lambda=\Lambda P_{t} \tag{1.1}
\end{equation*}
$$

in the sense of the composition of Markov kernels. Depending on non-degeneracy properties of $\Lambda$, such a relation is more or less strong. Especially when Markov semigroups are described by their generators, (1.1) is often replaced by

$$
\begin{equation*}
\bar{L} \Lambda=\Lambda L, \tag{1.2}
\end{equation*}
$$

where $\bar{L}$ and $L$ are respectively the generators of $\left(\bar{P}_{t}\right)_{t \geqslant 0}$ and $\left(P_{t}\right)_{t \geqslant 0}$. But then one has to be more careful with the meaning of generators (e.g. in the sense of martingale problems) and their domains, in particular the domains are transported via (1.2).

To be more useful from a probabilist point of view, it is convenient to convert (1.2) into a coupling between $\left(\bar{X}_{t}\right)_{t \geqslant 0}$ and $\left(X_{t}\right)_{t \geqslant 0}$, two Markov processes respectively associated to $\bar{L}$ and $L$ (called the dual and primal processes), so that the following relations hold for the conditional laws:

$$
\begin{equation*}
\forall t \geqslant 0, \quad \mathcal{L}\left(X_{t} \mid \bar{X}_{[0, t]}\right)=\Lambda\left(\bar{X}_{t}, \cdot\right) . \tag{1.3}
\end{equation*}
$$

In addition, one asks that $\left(\bar{X}_{t}\right)_{t \geqslant 0}$ can be constructed from $\left(X_{t}\right)_{t \geqslant 0}$ in an adapted way, meaning

$$
\begin{equation*}
\forall t \geqslant 0, \quad \mathcal{L}\left(\bar{X}_{[0, t]} \mid X\right)=\mathcal{L}\left(\bar{X}_{[0, t]} \mid X_{[0, t]}\right) . \tag{1.4}
\end{equation*}
$$

Yor was wondering about such couplings between some piecewise linear Markov processes and squared Bessel processes, in order to simplify his approach to certain properties of the former processes similar to those of the latter, see the end of the introduction of [23].

Such couplings are crucial for the constructions of strong stationary times, as explained by Diaconis and Fill [9] in a discrete time and finite setting. More precisely, in this situation $X$ is an ergodic Markov chain with invariant probability $\pi$ and $\bar{X}$ is a Markov chain absorbed in a unique point. A strong stationary time $\tau$ for $\left(X_{t}\right)_{t \geqslant 0}$ is a finite stopping time for $\left(X_{t}\right)_{t \geqslant 0}$ (and some independent randomness) such that $\tau$ and $X_{\tau}$ are independent and $X_{\tau}$ is distributed according to $\pi$. Taking into account (1.3) and (1.4), one can see that the absorption time for $\left(\bar{X}_{t}\right)_{t \geqslant 0}$ is a strong stationary time for $\left(X_{t}\right)_{t \geqslant 0}$.

Strong stationary times are important for two reasons (cf. Diaconis and Fill [9]):

- They enable to sample exactly the invariant probability $\pi$, contrary to the usual approximations provided by Monte Carlo techniques.
- They provide a probabilistic alternative to functional analysis approaches for the quantitative investigation of convergence to equilibrium. More precisely, for any strong stationary time $\tau$, we have

$$
\forall t \geqslant 0, \quad \mathfrak{s}\left(\mathcal{L}\left(X_{t}\right), \pi\right) \leqslant \mathbb{P}[\tau>t],
$$

where the separation discrepancy $\mathfrak{s}(\mu, \pi)$ between two probability measures $\mu$ and $\pi$ is defined by

$$
\mathfrak{s}(\mu, \pi):=\underset{\pi}{\operatorname{ess} \sup }\left(1-\frac{d \mu}{d \pi}\right)
$$

(where $d \mu / d \pi$ is the Radon-Nikodym density). The separation discrepancy dominates the total variation norm and gives positivity properties of $\mu$ with respect to $\pi$. In the context of convergence to equilibrium, it is very difficult to estimate the discrepancy of $\mathfrak{s}\left(\mathcal{L}\left(X_{t}\right), \pi\right)$ via functional inequality techniques (see e.g. the book [5] of Bakry, Gentil and Ledoux).

In the objective of constructing strong stationary times via intertwining duality, there are particular dual processes $\left(\bar{X}_{t}\right)_{t \geqslant 0}$ which are taking values in $\mathcal{M}$, the set of measurable subsets of $M$, but in general $\bar{M}$ is only a subset of $\mathcal{M}$, consisting in some regular subsets. The absorption set is the whole set $M$. The heuristic goal of intertwining duality is then to construct random subsets $\bar{X}_{t} \subset M$ such that $X_{t}$ is already at equilibrium in $\bar{X}_{t}$, for all $t \geqslant 0$, in such a way that $\bar{X}$ is itself Markovian and ends up covering the whole state space $M$.

In the diffusion context, set-valued intertwining dual processes started to be constructed in Fill and Lyzinski [11] and [15]. In [8], set-valued dual processes for diffusions on Riemannian manifolds were identified as stochastic perturbations of meancurvature flows. But the coupling of primal and dual processes were not considered in [8] and this is our present goal, mainly for Brownian motions on Riemannian manifolds. As we will see, there are numerous ways to construct such couplings (this is true in more general contexts, see [16] for the diversity of such couplings in a finite framework), but none of them is immediate and they are related to fine geometric features of the evolving subsets, such as their skeletons. We are thus to consider synchronous intertwined, free intertwined, mirror set-valued intertwined dual processes.

The reader must be warned that, as it stands now in the context of multidimensional diffusions, the set-valued dual processes are not defined up to the absorption time (except in symmetric settings), and as a consequence the same will be true for our couplings, which will be defined only up to some positive stopping times. We hope to investigate this point in future works, to end the construction of strong stationary times for Brownian motion on compact Riemannian manifolds, which remains our remote motivation. Other motivations for the couplings of primal and dual processes in the context of diffusions can be found in Machida [13] and [16].

Let us now present more precise definitions. Here the state space $M$ is a $d$-dimensional complete Riemannian manifold. Denote respectively by $\rho, \mu$ and $\underline{\mu}$, the Riemannian distance, the Lebesgue measure on $M$ and the corresponding ( $d-1$ )Hausdorff measure. The main objective of this paper is to construct couplings of primal diffusions processes with their set-valued dual intertwined processes. This will partially solve Conjecture 6 in [8] in the case of Brownian motion $\left(X_{t}\right)_{t \geqslant 0}$ and stochastic modified mean curvature flow $\left(D_{t}\right)_{t \geqslant 0}$ (which were generically denoted $\left(\bar{X}_{t}\right)_{t \geqslant 0}$ above). This conjecture says that an intertwined construction in the sense of Definition 1.1 is always possible.

Definition 1.1. - Consider a Markov process $D=\left(D_{t}\right)_{t \in[0, \tau]}$, with values in compact subsets of $M$ and continuous with respect to the Hausdorff topology, and where $\tau$ is an a.s. positive stopping time in the filtration $\mathscr{F}^{D}$ of $D$, serving as a lifetime for $D$. We say that a Brownian motion $X=\left(X_{t}\right)_{t \geqslant 0}$ in $M$ and $D$ are intertwined when 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 $X_{\tau^{\prime}} \in D_{\tau^{\prime}}$ ). More generally, for any $\mathscr{F}^{D}$-stopping time $\widetilde{\tau}$ smaller than $\tau$, we say that $X$ and $D$ are $\widetilde{\tau}$-intertwined when $X$ and $\left(D_{t}\right)_{t \in[0, \widetilde{\tau}]}$ are intertwined.

This is a generic definition, below stronger topologies on subsets of $M$ will be considered. Note that the above lifetime is not necessary the explosion time, i.e., the exit time from all compact sets for the considered topology. In the infinite dimensional state space of $D$, compactness does not seem an appropriate notion.

Also notice that $\widetilde{\tau}$-intertwining prevents $\left(X_{t}\right)_{t \geqslant 0}$ to have a lifetime smaller that $\widetilde{\tau}$. So we will never have to consider the lifetime of $\left(X_{t}\right)_{t \geqslant 0}$.

Our main results are Theorems 2.8, 2.12, 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 which is proved to be intertwined in Theorem 2.12, consists of the infinite-dimensional stochastic differential equation (2.6), 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 the precise requirements). Theorem 3.5 is obtained by specifying some functions $f$ approximating the distance to boundary. 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 (3.19)). Functions $f$ approximating the null function 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 of $\left(D_{t}\right)_{t \geqslant 0}$ is now independent from $\left(X_{t}\right)_{t \geqslant 0}$, while it is 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 [2].

To come back to our initial motivation, assume that $\left(X_{t}\right)_{t \geqslant 0}$ and $\left(D_{t}\right)_{t \geqslant 0}$ are intertwined, where the lifetime $\tau$ is the hitting/covering 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 (called the uniform distribution, which is invariant and reversible for the Brownian motion $\left.\left(X_{t}\right)_{t \geqslant 0}\right)$ and $\tau$ is a strong stationary time for $\left(X_{t}\right)_{t \geqslant 0}$. 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. These estimates will need geometric ingredients such as Ricci bounds and it will be interesting to see how they will enter the game.

The needs for couplings between primal and dual processes of a Markovian intertwining relation is illustrated by [3], where strong stationary times $\tau_{n}$ are constructed for the $n$-dimensional sphere (when the subset-valued dual is starting from a singleton), satisfying

$$
\mathbb{E}\left[\tau_{n}\right] \sim \frac{\ln (n)}{n}
$$

and for any $r>0$,

$$
\lim _{n \rightarrow \infty} \mathbb{P}\left[\tau_{n}>(1+r) \frac{\ln (n)}{n}\right]=\lim _{n \rightarrow \infty} \mathbb{P}\left[\tau_{n}<(1-r) \frac{\ln (n)}{n}\right]=0
$$

## 2. Intertwined dual processes: existence in connection with Stokes' formula

In this section we make a construction of intertwined processes $\left(X_{t}\right)_{t \geqslant 0}$ and $\left(D_{t}\right)_{t \geqslant 0}$ based on the Stokes' Formula (2.1) below. Consider a compact domain $D$ in $M$ with $C^{2}$ boundary. Let $f: D \rightarrow \mathbb{R}$ a $C^{2}$ function such that $\left.\nabla f\right|_{\partial D}=N^{D}$ the normal inward vector on boundary. Then by Stokes' formula, for any $C^{2}$ function $g: D \rightarrow \mathbb{R}$,

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

For $\alpha \in(0,1)$, denote by $\mathscr{D}^{2+\alpha}$ the set of compact connected subsets $D$ of $M$ with $C^{2+\alpha}$ boundary. It will be more convenient to work with this state space (endowed with its natural topology) than with the larger one considered in Definition 1.1. Let us even restrict it further:

We fix a point $o \in M$ for convenience.

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

- $D \subset B(o, 1 / \varepsilon)$ the Riemannian ball centered at $o$ with radius $1 / \varepsilon$;
- $\rho(\partial D, S(D)) \geqslant \varepsilon$, 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 $(D)^{c}$;
- the coefficients of the $\alpha$-Hölderianity of the second fundamental form of $\partial D$ are bounded by $1 / \varepsilon$.
The set $\mathcal{F}^{\alpha, \varepsilon}$ will serve as the state space of the set-valued process $\left(\widetilde{D}_{t}\right)_{t \in\left[0, \tau_{\varepsilon}\right]}$ and $\tau_{\varepsilon} \in(0,+\infty]$ will be the exiting time from $\mathcal{F}^{\alpha, \varepsilon}$. This process will be a diffusion, i.e., a Markov process with continuous trajectories (for the topology inherited from $\mathscr{D}^{2+\alpha}$ ), and its generator $\widetilde{\mathscr{L}}$ will be defined later in (2.8). 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:

$$
\begin{equation*}
\mathcal{V}_{\delta}^{2+\beta}\left(D_{0}\right):=\left\{\operatorname{int}\left(\exp _{\partial D_{0}}(f)\right), f \in C^{2+\beta}\left(\partial D_{0}\right),\|f\|_{C^{2+\beta}\left(\partial D_{0}\right)}<\delta\right\}, \tag{2.2}
\end{equation*}
$$

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 relatively compact open subset of $M$ with boundary $\exp _{\partial D_{0}}(f)$.

Let $\eta\left(\partial D_{0}\right)>0$ be the maximal radius for a tubular neighborhood of $\partial D_{0}$ on which the signed distance to $\partial D_{0}$ is regular. Alternatively, $\eta\left(\partial D_{0}\right)$ is the distance between $\partial D_{0}$ and the union of inner and outer skeleton of $D_{0}$. Notice that $\delta<\eta\left(\partial D_{0}\right)$ guarantees 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 $\eta(\partial D) \geqslant \varepsilon$. The map which to $\left\{f \in C^{2+\beta}\left(\partial D_{0}\right),\|f\|_{C^{2+\beta}\left(\partial D_{0}\right)}<\delta\right\}$ associates $D \in \mathcal{V}_{\delta}^{2+\beta}\left(D_{0}\right)$ via (2.2) is one-to-one since for such a $D$, the corresponding function $f$ is characterized by the fact that at a point $z$ of $\partial D$ which projects onto $\pi(z) \in \partial D_{0}, f(\pi(z))$ is the signed distance from $\partial D_{0}$ to $z$, positive when $z \in D_{0}$. This is a particular case of [7, Th. 1.5].

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

$$
\begin{equation*}
d_{\beta, D_{0}}\left(D_{1}, D_{2}\right):=\left\|f_{1}-f_{2}\right\|_{C^{2+\beta}\left(\partial D_{0}\right)} . \tag{2.3}
\end{equation*}
$$

## Assumption 2.2

- The function

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

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 defined by the distances $d_{\alpha, D}$ ). The functions $f^{D}$ satisfy

$$
\left\|\nabla f^{D}\right\|_{\infty} \leqslant 1
$$

and coincide with the signed distance to the boundary $\rho_{\partial D}^{+}$(positive inside $D$ and negative outside) 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 $\operatorname{Hess} f^{D}$ are uniformly bounded over $\mathcal{F}^{\alpha, \varepsilon}$.

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

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

where $\Gamma\left(T M \otimes\left(\mathbb{R}^{m}\right)^{*}\right)$ denotes 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

$$
\begin{aligned}
\sigma^{D}(x): \mathbb{R} \times \mathbb{R}^{m} & \longrightarrow T_{x} M \\
\left(w_{0}, w\right) & \longmapsto w_{0} \nabla f^{D}(x)+\sigma_{\mathrm{c}}^{D}(x)(w)
\end{aligned}
$$

satisfies

$$
\forall x \in D, \quad \sigma^{D}\left(\sigma^{D}\right)^{*}(x)=\operatorname{Id}_{T_{x} M}
$$

Remark 2.3. - The first condition of Assumption 2.2 implies that

$$
\begin{align*}
\left.\nabla f^{D}\right|_{\partial D} & =\left.\left(\nabla \rho_{\partial D}^{+}\right)\right|_{\partial D}\left(=N^{D}\right) \quad \text { and } \\
\left.\Delta f^{D}\right|_{\partial D} & =\left.\left(\Delta \rho_{\partial D}^{+}\right)\right|_{\partial D}\left(=-h^{D}\right), \tag{2.4}
\end{align*}
$$

where $h^{D}$ stands for the mean curvature on $\partial D$ : at $x \in \partial D, h^{D}(x)$ is the trace of the second fundamental form of $\partial D$, it can alternatively be described as the sum of the principal curvatures in 2-planes directions containing $N^{D}(x)$. The sign convention is that $h^{D}>0$ when $D$ is convex. 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:

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

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 $u \in T_{x} M$,

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

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 $\left\|N^{D}(x)\right\|=1$ :

$$
\begin{equation*}
0=\left\langle\nabla f^{D}(x), \sigma_{c}^{D}(x)\left(e_{i}\right)\right\rangle, \quad i=1, \ldots m \tag{2.5}
\end{equation*}
$$

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 $D_{\varepsilon}:=\{x \in M, \rho(x, \partial D)<\varepsilon\}$. Let $h_{\varepsilon}$ be an odd smooth nondecreasing function from $\mathbb{R}$ 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 [4, Prop. 3.2], where $m+1, \nabla f^{D},\left(\sigma_{c}^{D}\left(e_{1}\right), \ldots, \sigma_{c}^{D}\left(e_{m}\right)\right)$ here is denoted $m, \sigma_{1},\left(\sigma_{2}, \ldots, \sigma_{m}\right)$ there. The wanted regularity in $D$ is easily checked.

Let $\left(W_{t}\right)_{t \geqslant 0}$ and $\left(W_{t}^{m}\right)_{t \geqslant 0}$ 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 $y \in \partial D_{t}$ :

$$
\left\{\begin{align*}
d X_{t} & =\left(\nabla f^{D_{t}}\left(X_{t}\right) d W_{t}+\sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right)  \tag{2.6}\\
d \partial D_{t}(y) & =N^{D_{t}}(y)\left(d W_{t}+\left(\frac{1}{2} h^{D_{t}}(y)+\Delta f^{D_{t}}\left(X_{t}\right)\right) d t\right)
\end{align*}\right.
$$

started at a 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 $D_{0}$. This assumption is essential and will be made all along the paper. The notation $d \partial D_{t}(y)=\left(d \partial D_{t}\right)(y)$ stands for an infinitesimal move of the boundary $\partial D_{t}$ at point $y$ and is rigorously presented in Appendix B, see (B.3). The second equation in (2.6) and (2.7) below are stochastic differential equations in $\mathscr{D}^{2+\alpha}$, and a geometric way to represent stochastic partial differential equations locally defined in $C^{\alpha, 2+\alpha}\left([0, \infty) \times \partial D_{0}\right)$. Similar equations can be found in [12, App. A.2].

In fact, as in Definition 2.1, the evolution equation (2.6) is implicitly considered only up to the a.s. positive 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 (2.6), 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 a solution to the evolution equation

$$
\begin{equation*}
\forall t \leqslant \widetilde{\tau}_{\varepsilon}, \forall y \in \partial \widetilde{D}_{t}, \quad d \partial \widetilde{D}_{t}(y)=N^{\widetilde{D}_{t}}(y)\left(d \widetilde{W}_{t}+\left(\frac{1}{2} h^{\widetilde{D}_{t}}(y)-\frac{\mu^{\partial \widetilde{D}_{t}}\left(\partial \widetilde{D}_{t}\right)}{\mu\left(\widetilde{D}_{t}\right)}\right) d t\right), \tag{2.7}
\end{equation*}
$$

where $\left(\widetilde{W}_{t}\right)_{t \geqslant 0}$ is a real-valued Brownian motion and where $\widetilde{\tau}_{\varepsilon}$ is the exit time from $\mathcal{F}^{\alpha, \varepsilon}$.

Notice that the equation for $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ does no longer depend of $\left(X_{t}\right)_{t \geqslant 0}$, so if the solution is unique, $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ will be Markovian. It is [8, Eq. (44)] (up to a time scaling by 2). Theorem 40 of [8] (where (44) has been rewritten as (79)) proves local existence of a solution. The second term in the right of (2.7) is the one of mean curvature flow (with $1 / 2$ in front of it). The first term is a stochastic perturbation, uniform in the normal direction. The equation for stochastic front propagation in [12] has exactly these two terms. The last term in our equation is also uniform in the normal direction. It can be seen as a conditioning which prevents the solution to implode. One of the main goals of this article will be to prove that the solution to the second equation in (2.6) has same law as the solution to (2.7).

Theorem 2.5. - Fix $\alpha \in(0,1)$ and $\varepsilon>0$. Then (2.7) 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[8, Th. 22]. 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

$$
\begin{align*}
\widetilde{\mathscr{L}}\left[F_{k}\right](D) & :=\underline{\mu}^{\partial D}(k) \frac{\underline{\mu}^{\partial D}(\partial D)}{\mu(D)}-\frac{1}{2} \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right),  \tag{2.8}\\
\Gamma_{\widetilde{\mathscr{L}}}\left[F_{k}, F_{g}\right](D) & :=\int_{\partial D} k d \underline{\mu} \int_{\partial D} g d \underline{\mu} . \tag{2.9}
\end{align*}
$$

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

$$
\begin{equation*}
\widetilde{\mathscr{L}}[\mathfrak{F}]:=\sum_{l=1}^{n} \partial_{j} \mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right) \widetilde{\mathscr{L}}\left[F_{k_{l}}\right]+\sum_{j, l \in \llbracket 1, n \rrbracket}^{n} \partial_{j, l} \mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right) \Gamma_{\mathscr{\mathscr { L }}}\left[F_{k_{j}}, F_{k_{l}}\right] . \tag{2.10}
\end{equation*}
$$

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

$$
\begin{equation*}
\Gamma_{\mathscr{L}}[\mathfrak{F}, \mathfrak{G}]:=\sum_{l \in \llbracket n \rrbracket, j \in \llbracket m \rrbracket} \partial_{l} \mathfrak{f}\left(F_{k_{1}}, \ldots, F_{k_{n}}\right) \partial_{j} \mathfrak{g}\left(F_{g_{1}}, \ldots, F_{g_{m}}\right) \Gamma_{\mathscr{L}}\left[F_{k_{l}}, F_{g_{j}}\right] . \tag{2.11}
\end{equation*}
$$

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 [8, Rem. 2]. More generally, the forms of (2.10) and (2.11) 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 [5].

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 $\tau_{\varepsilon}$ (which was defined after Definition 2.1) or $\widetilde{\tau}_{\varepsilon}$ (defined after (2.7)).

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 (2.6) has a solution $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ started at $D_{0} \in \mathscr{F}^{\alpha, \varepsilon}, X_{0} \sim \mathscr{U}\left(D_{0}\right)$.

Proof. - 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

$$
\left\{\begin{align*}
d \partial D_{t}(y) & =N^{D_{t}}(y)\left(d \widehat{W}_{t}+\left(\frac{1}{2} h^{D_{t}}(y)-\frac{\underline{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}\right) d t\right)  \tag{2.12}\\
d X_{t} & =\left(\nabla f^{D_{t}}\left(X_{t}\right)\left[d \widehat{W}_{t}-\left(\frac{\underline{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}+\Delta f^{D_{t}}\left(X_{t}\right)\right) d t\right]\right. \\
& \left.+\sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right)
\end{align*}\right.
$$

for $\left(\widehat{W}_{t}\right)_{t \geqslant 0}$ and $\left(W_{t}^{m}\right)_{t \geqslant 0}$ independent Brownian motions. Notice that the first equation is the same as (2.7). Thus due to Theorem 2.5, $\left(D_{t}\right)_{t \geqslant 0}$ is a diffusion process with generator $\widetilde{\mathscr{L}}$. Then given $\left(D_{t}\right)_{t \geqslant 0}$, the equation for $\left(X_{t}\right)_{t \geqslant 0}$

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

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 (2.5) which yields on boundary $\left\langle N^{D_{t}}\left(X_{t}\right), \sigma_{\mathrm{c}}^{D_{t}}\left(X_{t}\right) d W_{t}^{m}\right\rangle=0$,

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

where we used (2.12) and (2.4). We also have no covariation since the martingale part of $d \partial D_{t}$ acts on the normal flow only, and any normal flow

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

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 (2.12), make by Girsanov theorem a change of probability such that $\left(W_{t}, W_{t}^{m}\right)_{t \geqslant 0}$ is a Brownian motion where

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

We get a solution to (2.6) in the new probability.

Proposition 2.9. - Let $\left(D_{t}\right)_{t \geqslant 0}$ satisfy

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

for some Brownian motion $\left(W_{t}\right)_{t \geqslant 0}$ and some adapted locally bounded real-valued process $b_{t}$. Let $\mu_{t}=\mu^{D_{t}}$ be the Lebesgue measure on $D_{t}$ and $\bar{\mu}_{t}=\bar{\mu}^{D_{t}}=\mathscr{U}\left(D_{t}\right)=$ $\mu^{D_{t}} / \mu\left(D_{t}\right)$. Denote by $\mu_{t}=\mu^{\partial D_{t}}$ the Lebesgue measure on $\partial D_{t}$ and $\bar{\mu}_{t}=\bar{\mu}^{\partial D_{t}}=$ $\mu^{\partial D_{t}} / \mu\left(D_{t}\right)$. Let $k$ be a smooth function of $M$. Then

$$
\begin{equation*}
d \mu_{t}(k)=-\underline{\mu}_{t}(k) d W_{t}-\frac{1}{2}\left(2 b_{t} \underline{\mu}_{t}(k)+\underline{\mu}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right)\right) d t \tag{2.13}
\end{equation*}
$$

and

$$
\begin{align*}
& d \bar{\mu}_{t}(k)=\left(-\overline{\underline{\mu}}_{t}(k)+\bar{\mu}_{t}(k) \overline{\underline{\mu}}_{t}\left(\partial D_{t}\right)\right) d W_{t}-\frac{1}{2} \overline{\underline{\mu}}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right) d t  \tag{2.14}\\
&+\left(\bar{\mu}_{t}\left(\partial D_{t}\right)+b_{t}\right)\left(-\overline{\underline{\mu}}_{t}(k)+\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right)\right) d t
\end{align*}
$$

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

$$
\begin{equation*}
d \bar{\mu}_{t}(k)=\left(-\overline{\underline{\mu}}_{t}(k)+\bar{\mu}_{t}(k) \bar{\mu}_{t}\left(\partial D_{t}\right)\right) d W_{t}-\frac{1}{2} \overline{\underline{\mu}}_{t}\left(\left\langle d k, N^{D_{t}}\right\rangle\right) d t \tag{2.15}
\end{equation*}
$$

Proof. - Let us first work at fixed time $t \geqslant 0$. Denote $D=D_{t}$ and adopt the corresponding notations presented in Appendix A. For $k$ 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

$$
F(r, k)=\int_{D(r)} k d \mu
$$

We have

$$
F(r, k)=\int_{\partial D}\left(\int_{r}^{\tau(y)} k(\psi(s)(y)) e^{-\int_{0}^{s} h^{D}(\psi(u)(y)) d u} d s\right) \underline{\mu}(d y)
$$

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). The mapping $h^{D}$ is defined in (A.7) and is an extension of the mean curvature on the boundary $\partial D$ : it corresponds to the mean curvature for the foliation induced by the $\partial D(r), r \in \mathbb{R}$ sufficiently small. With this formulation we can differentiate with respect to $r$, to obtain

$$
F^{\prime}(r, k)=-\int_{\partial D} k(\psi(r, y)) e^{-\int_{0}^{r} h^{D}(\psi(s)(y) d s} \underline{\mu}(d y)
$$

Differentiating again we get

$$
F^{\prime \prime}(r, k)=-\int_{\partial D}\left(\left\langle d k, \partial_{r} \psi(r, y)\right\rangle-(k h)(\psi(r, y))\right) e^{-\int_{0}^{r} h^{D}(\psi(s)(y)) d s} \underline{\mu}(d y)
$$

In particular,

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

This allows us to compute

$$
d\left(F\left(W_{t}, k\right)\right)=F^{\prime}\left(W_{t}, k\right) d W_{t}+\frac{1}{2} F^{\prime \prime}\left(W_{t}, k\right) d t
$$

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

$$
d \mu_{t}(k)=\int_{\partial D_{t}}-k(y)\left\langle d \partial D_{t}(y), N^{D_{t}}(y)\right\rangle+\frac{1}{2}\left(k h^{D_{t}}-\left\langle d k, N^{D_{t}}\right\rangle\right)(y) \mu_{t}(d y) .
$$

This yields

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

which gives (2.13). In particular, taking $k \equiv 1$ we obtain

$$
\begin{equation*}
d \mu\left(D_{t}\right)=\mu_{t}\left(\partial D_{t}\right)\left(-d W_{t}-b_{t} d t\right) \tag{2.16}
\end{equation*}
$$

Now we can compute

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

This yields (2.14).
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 (2.6) 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

$$
F_{k}(D):=\int_{D} k d \mu,
$$

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

$$
\begin{align*}
\mathscr{L}\left(g F_{k}\right)(x, D)=-g(x) \Delta f^{D}(x) \underline{\mu}^{\partial D}(k) & -\frac{1}{2} g(x) \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)  \tag{2.17}\\
& +\frac{1}{2} F_{k}(D) \Delta g(x)-\underline{\mu}^{\partial D}(k)\left\langle\nabla g, \nabla f^{D}\right\rangle(x) .
\end{align*}
$$

Proof. - From (2.6) and (2.13) with $b_{t}=\Delta f^{D_{t}}\left(X_{t}\right)$ we have

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

This implies that

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

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

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

Consequently, using

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

we get (2.17).
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 (2.17) 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 (2.6) 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 (2.7)), $\widetilde{\nu}_{t}$ its law at time $t$ and let

$$
\nu_{t}(d D, d x):=\widetilde{\nu}_{t}(d D) \mathscr{U}(D)(d x) .
$$

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

$$
\begin{equation*}
\partial_{t} \nu_{t}\left(g F_{k}\right)=\nu_{t}\left(\mathscr{L}\left(g F_{k}\right)\right) \tag{2.18}
\end{equation*}
$$

Proof. - Integrating (2.17) in $x$ with respect to the uniform law $\bar{\mu}^{D}:=\mathscr{U}(D)$ in $D$ yields

$$
\begin{align*}
&-\bar{\mu}^{D}\left(g \Delta f^{D}\right) \underline{\mu}^{\partial D}(k)-\frac{1}{2} \bar{\mu}^{D}(g) \underline{\mu}^{\partial D}\left(\left\langle\nabla k, N^{D}\right\rangle\right)  \tag{2.19}\\
&+\frac{1}{2} F_{k}(D) \bar{\mu}^{D}(\Delta g)-\underline{\mu}^{\partial D}(k) \bar{\mu}^{D}\left(\left\langle\nabla g, \nabla f^{D}\right\rangle\right)
\end{align*}
$$

By Stokes theorem,

$$
\bar{\mu}^{D}\left(g \Delta f^{D}+\left\langle\nabla g, \nabla f^{D}\right\rangle\right)=\bar{\mu}^{\partial D}\left(g\left\langle\nabla f^{D},-N^{D}\right\rangle\right)=-\bar{\mu}^{\partial D}(g),
$$

so the expression (2.19) writes

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

On the other hand,

$$
\nu_{t}\left(g F_{k}\right)=\widetilde{\nu}_{t}\left[\bar{\mu}^{D_{t}}[g] F_{k}\right],
$$

which implies that

$$
\begin{equation*}
\partial_{t} \nu_{t}\left(g F_{k}\right)=\partial_{t} \widetilde{\nu}_{t}\left(\left(\bar{\mu}^{D_{t}}(g) F_{k}\right)=\widetilde{\nu}_{t}\left(\widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g) F_{k}\right)\right) .\right. \tag{2.20}
\end{equation*}
$$

By (2.15),

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

so, taking into account (2.9),

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

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

$$
H\left(D_{t}\right)=\widetilde{\mathscr{L}}\left(\bar{\mu}^{D_{t}}(g) F_{k}\right),
$$

which together with (2.20) proves (2.18).
Theorem 2.12. - Let $(x, D) \mapsto f^{D}(x)$ and $(x, D) \mapsto \sigma_{\mathrm{c}}^{D}(x)$ satisfy Assumption 2.2. Consider a solution $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ to equation (2.6) started at $D_{0} \in \mathscr{F}{ }^{\alpha, \varepsilon}, X_{0} \sim \mathscr{U}\left(D_{0}\right)$. Then for all $t \geqslant 0,\left(D_{t}, X_{t}\right)$ 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 $\left(D_{t}\right)_{t \geqslant 0}$ is a diffusion with generator $\widetilde{\mathscr{L}}$.

Proof. - 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

$$
G(g, k, t)(s)=\nu_{s}\left(\mathscr{P}_{t-s}\left(g F_{k}\right)\right), \quad s \in[0, t] .
$$

We compute

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

where we used Proposition 2.11 in the first term of the second line, and 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

$$
\nu_{0} \mathscr{P}_{t}\left(g F_{k}\right)=\nu_{t}\left(g F_{k}\right),
$$

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 $\left(D_{t}\right)_{t \geqslant 0}$ 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 $\left(\widetilde{W}_{t}\right)_{t \geqslant 0}$ be a real-valued Brownian motion and $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ be the solution of (2.7) started at $\widetilde{D}_{0}$, with driving Brownian motion $\left(\widetilde{W}_{t}\right)_{t \geqslant 0}$.

Assumption 3.1. - Fix $\alpha \in(0,1)$ and $\varepsilon>0$. 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 continuous 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, Brownian motions with probability one never hit the singular part of $S\left(\widetilde{D}_{t}\right)$.

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

Notice that [1, Th. 1.1] proves the first part of the conjecture, i.e., the continuity of $D \mapsto S(D)$, in the case where $M=\mathbb{R}^{d}$ endowed with a possibly varying Riemannian metric. All examples in Section 5 together with the study of the motion of the skeleton in Appendix B make us believe that Conjecture 3.2 is true. In particular, Section 5.4 provides a large class of examples in $\mathbb{R}^{2}$ which do not reduce to finite dimensional processes, some of them having infinite lifetime. They are characterized by the fact that the motion of skeleton can be explicitly described. The considered skeletons have sufficient number of symmetries. For simplicity we considered $n$-branches skeletons, but we could consider trees with as many ramifications as we want. We could also replace $\mathbb{R}^{2}$ by the hyperbolic plane or the two dimensional sphere, as well as dimension 2 by higher dimension. All these situations would furnish true infinite dimensional set-valued processes, some of them with completely describable skeleton.

However a better knowledge of skeletons is necessary to solve the conjecture in the general situation. We believe that the process $\left(S\left(\widetilde{D}_{t}\right)\right)_{t \geqslant 0}$ takes its values in a set of regular stratified spaces, and that it has absolutely continuous variation in this space.

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 $\mathbb{R}^{d} \ni u \mapsto \varphi_{\delta}(|u|)$ is smooth and $\int_{\mathbb{R}^{d}} \varphi_{\delta}(|u|) d u=1$ (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 $T_{x} M$ ).
Define

$$
\left.e(\delta):=\sup \{\| \|(\nabla \exp )(u)] \|, x \in B(o, 1 / \varepsilon), u \in B_{x}(0, \delta) \subset T_{x} M\right\}
$$

where $\nabla \exp (u): T_{x} M \rightarrow T_{\exp _{x}(u)} M$ is the covariant derivative of exp with respect to the base point, $\|\|\cdot\|\|$ is the operator norm, when $T_{x} M$ and $T_{\exp _{x}(u)} M$ are endowed with their Euclidean structures, and $B_{x}(0, \delta)$ is the open ball in $T_{x} M$ with center 0 and radius $\delta$. Recall 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,

$$
0 \leqslant e(\delta)-1 \leqslant C_{1}^{\prime} \delta,\left|c_{\delta}-1\right| \leqslant C_{1}^{\prime \prime} \delta .
$$

Proof. - The inequalities of the first line are well-known 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 $B(o, 1 / \varepsilon) \times B(o, 1 / \varepsilon)$. Then, for any $x \in M$, we can rewrite (3.1) under the form

$$
\begin{aligned}
f_{\delta}(x, D)=\ell_{\varepsilon}(x, D) \rho_{\delta}( & x, \partial D) \\
& +\left(1-\ell_{\varepsilon}(x, D)\right) \int_{M} \varphi_{\delta}\left(\left|\exp _{x}^{-1}(y)\right|\right) \rho_{\delta}(y, \partial D) J \exp _{x}^{-1}(y) d y
\end{aligned}
$$

where $J \exp _{x}^{-1}$ is the absolute value of the determinant of the Jacobian of $\exp _{x}^{-1}(\cdot)$.
The interest of all these preparations is:

## Proposition 3.4

For all $\delta>0$ sufficiently small, 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

$$
\begin{equation*}
\left|f_{\delta}(x, D)-\rho(x, \partial D)\right| \leqslant C_{1} \delta \tag{3.3}
\end{equation*}
$$

- the differential 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 $\partial D$ :

$$
\begin{equation*}
\left\langle d_{2} f_{\delta}(x, D), K\right\rangle \leqslant C_{4}\|K\|_{\infty} \quad \text { and } \quad\left\|\nabla_{2} d_{2} f_{\delta}(x, D)\left(N_{\partial D}, N_{\partial D}\right)\right\| \leqslant C_{4} \tag{3.4}
\end{equation*}
$$

for a $C_{4}$ not depending on $x, D, \delta$. The second term is the second derivative along the inward normal flow on $D$.

Proof. - We first prove $\left\|d_{1} f_{\delta}(x, D)\right\| \leqslant 1, d_{1}$ denoting the differential with respect to the first or the $x$ variable. For $x \in B(o, 1 / \varepsilon)$ we have

$$
\begin{aligned}
& d_{1} f_{\delta}(x, D)= \ell_{\varepsilon}(x, D) \\
&+\left(1-\ell_{1} \rho_{\delta}(x, \partial D)\right. \\
&+(x, D)) d_{1}\left(\int_{T_{x} M} \varphi_{\delta}(|u|) \rho_{\delta}\left(\exp _{x}(u), \partial D\right) d u\right) \\
&+d_{1} \ell_{\varepsilon}(x, D) \int_{T_{x} M} \varphi_{\delta}(|u|)\left(\rho_{\delta}(x, \partial D)-\rho_{\delta}\left(\exp _{x}(u), \partial D\right)\right) d u .
\end{aligned}
$$

Notice that if $x^{\prime}$ is close to $x$ and $v_{x, x^{\prime}}: T_{x} M \rightarrow T_{x^{\prime}} M$ is the parallel transport along the minimal geodesic from $x$ to $x^{\prime}$, then

$$
\int_{T_{x^{\prime}} M} \varphi_{\delta}(|u|) \rho_{\delta}\left(\exp _{x^{\prime}}(u), \partial D\right) d u=\int_{T_{x} M} \varphi_{\delta}(|u|) \rho_{\delta}\left(\exp _{x^{\prime}}\left(l_{x, x^{\prime}}(u), \partial D\right) d u\right.
$$

Taking the differential with respect to $x^{\prime}$ at $x^{\prime}=x$ and using $\left.\nabla_{x^{\prime}}\right|_{x^{\prime}=x} \imath_{x, x^{\prime}}=0$ by definition of parallel transport yields

$$
d_{1}\left(\int_{T_{x} M} \varphi_{\delta}(|u|) \rho_{\delta}\left(\exp _{x}(u), \partial D\right) d u\right)=\int_{T_{x} M} \varphi_{\delta}(|u|) d_{1} \rho_{\delta}((\nabla \exp )(u), \partial D) d u
$$

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

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

If $\rho(x, \partial D) \geqslant \varepsilon / 2$ then for $\delta \leqslant \varepsilon / 8$, we have $\rho\left(\exp _{x}(u), \partial D\right) \geqslant 3 \varepsilon / 8$ for $u \in T_{x} M$ with $|u| \leqslant \delta$. It follows

$$
\leqslant 1 .
$$

It is easily checked that the function $f_{\delta}$ satisfies the other properties of Assumption 2.2. Let us check that it also satisfies (3.3).

$$
\begin{aligned}
& \left\|d_{1} f_{\delta}(x, D)\right\| \leqslant \ell_{\varepsilon}(x) e^{-1}(\delta)\left(1-\delta\left\|d_{1} \ell_{\varepsilon}\right\|_{\infty}\right) \\
& +\left(1-\ell_{\varepsilon}(x)\right) \int_{T_{x} M} \varphi_{\delta}(|u|) c_{d}\|(\nabla \exp )(u)\| d u \\
& +\left\|d_{1} \ell_{\varepsilon}(x)\right\|_{\infty} \int_{T_{x} M} \varphi_{\delta}(|u|) \delta d u
\end{aligned}
$$

We have

$$
\begin{align*}
& f_{\delta}(x, D)-\rho_{\delta}(x, \partial D)  \tag{3.5}\\
& \quad=\left(1-\ell_{\varepsilon}(x, D)\right) \int_{T_{x} M} \varphi_{\delta}(|u|)\left(\rho_{\delta}\left(\exp _{x}(u), \partial D\right)-\rho_{\delta}(x, \partial D)\right) d u,
\end{align*}
$$

which implies

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

On the other hand,

$$
\left|\rho(x, \partial D)-\rho_{\delta}(x, \partial D)\right| \leqslant\left(1-c_{\delta}\right) \max (2 / \varepsilon, 3 \varepsilon / 8) \leqslant C_{1}^{\prime \prime \prime} \delta
$$

for some constant $C_{1}^{\prime \prime \prime}>0$ (depending on $\varepsilon$ ). This yields (3.3) with $C_{1}:=1+C_{1}^{\prime \prime \prime}$.
For proving (3.4), we take a vector field $K(y)=k(y) N(y), y \in \partial D$ and compute

$$
\left\langle d_{2} \rho(x, \partial D), K\right\rangle=\langle-N(P(x)), K(P(x))\rangle=-k(P(x)),
$$

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

$$
\nabla_{2} d_{2} \rho(x, \partial D)\left(N_{\partial D}, N_{\partial D}\right)=0
$$

Remarking that $\left\|d_{2} \ell_{\varepsilon}(x, D)\right\|$ is bounded by $\left\|h_{\varepsilon}^{\prime}\right\|_{\infty}$, we get (3.4) via a straightforward computation.

Theorem 3.5. - Fix $D_{0}=\widetilde{D}_{0} \in \widetilde{\mathcal{F}}^{\alpha, \varepsilon}$ 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

$$
\begin{align*}
d \partial D_{t}(y)=N^{D_{t}}(y)\left(\left\langle d X_{t}, N^{D_{t}}\left(X_{t}\right)\right\rangle+\left(\frac{1}{2} h^{D_{t}}(y)-\right.\right. & \left.h^{D_{t}}\left(X_{t}\right) \mathbb{1}_{D_{t} \backslash S_{t}}\left(X_{t}\right)\right) d t  \tag{3.6}\\
& \left.-2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X)\right) .
\end{align*}
$$

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 $x \in S_{t}$ (recall $S_{t}$ is the regular skeleton of $D_{t}$, see Appendix A). 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 $S_{t}:=S\left(D_{t}\right)$ :

$$
\begin{equation*}
L_{t}^{S_{t}}(X)=\lim _{\beta \searrow 0} \frac{1}{2 \beta} \int_{0}^{t} 1_{\left\{X_{s} \in S_{s}^{\beta}\right\}} d s \tag{3.7}
\end{equation*}
$$

$S_{s}^{\beta}$ 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 $S^{D}$

$$
\nabla \rho_{\partial D}(x)=N^{D}(x) \quad \text { and } \quad \Delta \rho_{\partial D}(x)=-h^{D}(x)
$$

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

$$
" \Delta \rho_{\partial D_{t}}\left(X_{t}\right) d t "=-2 \sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X) .
$$

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 (2.6)

$$
\begin{equation*}
d \partial D_{t}^{\delta}(y)=N^{D_{t}^{\delta}}(y)\left(d W_{t}^{\delta}+\left(\frac{1}{2} h^{D_{t}^{\delta}}(y)+\Delta f_{\delta}^{D_{t}^{\delta}}\left(X_{t}^{\delta}\right)\right) d t\right) \tag{3.8}
\end{equation*}
$$

for some Brownian motion $W_{t}^{\delta}$. On the other hand, from Proposition 2.11 and (2.1), $\left(D_{t}^{\delta}\right)_{t \geqslant 0}$ satisfies equation (2.7):

$$
d \partial D_{t}^{\delta}(y)=N^{D_{t}^{\delta}}(y)\left(d \widetilde{W}_{t}^{\delta}+\left(\frac{1}{2} h^{D_{t}^{\delta}}(y)-\frac{\mu^{\partial D_{t}^{\delta}}\left(\partial D_{t}^{\delta}\right)}{\mu\left(D_{t}^{\delta}\right)}\right) d t\right)
$$

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

$$
\begin{equation*}
d \widetilde{W}_{t}^{\delta}=d W_{t}^{\delta}+\Delta f_{\delta}^{D_{t}^{\delta}}\left(X_{t}\right) d t+\frac{\mu^{\partial D_{t}^{\delta}}\left(\partial D_{t}^{\delta}\right)}{\mu\left(D_{t}^{\delta}\right)} d t \tag{3.9}
\end{equation*}
$$

A remarkable fact about all $\left(X_{t}^{\delta}, D_{t}^{\delta}\right)_{t \geqslant 0}$ is that their marginals are constant in law: for the second marginal we use Proposition F. 2 which states that the martingale problem associated to $\breve{\mathscr{L}}$ is well posed, and this implies uniqueness 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

$$
\begin{equation*}
\left(\left(X_{t}^{\delta}, D_{t}^{\delta}, W_{t}^{\delta}, \widetilde{W}_{t}^{\delta}, W_{t}^{\delta, m}\right)_{t \geqslant 0}, \tau_{\varepsilon}^{\delta}\right) \tag{3.10}
\end{equation*}
$$

is tight (in (3.10) the Brownian motions $\left(W_{t}^{\delta}\right)_{t \geqslant 0}$ and $\left(W_{t}^{\delta, m}\right)_{t \geqslant 0}$ are the ones defined by equation (2.6)). Denote by

$$
\begin{equation*}
\left(\left(X_{t}, D_{t}, W_{t}, \widetilde{W}_{t}, W_{t}^{m}\right)_{t \geqslant 0}, \tau_{\varepsilon}\right) \tag{3.11}
\end{equation*}
$$

a limiting point. Let us prove the intertwining.
Using Proposition 2.11, for any smooth functions $g$ and $k$ on $M$, for any $t \geqslant 0$,

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

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

$$
d \partial D_{t}(y)=N^{D_{t}}(y)\left(d \widetilde{W}_{t}+\left(\frac{1}{2} h^{D_{t}}(y)-\frac{\mu^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)}\right) d t\right)
$$

We need to work with real-valued processes: we have from (2.16), for all $\delta>0$,

$$
\begin{equation*}
\int_{0}^{t} \frac{d \mu\left(D_{s}^{\delta}\right)}{\underline{\mu}\left(\partial D_{s}^{\delta}\right)}=-W_{t}^{\delta}-\int_{0}^{t} \Delta_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right) d s \tag{3.12}
\end{equation*}
$$

This together with (3.9) yields

$$
\begin{equation*}
d \partial D_{t}^{\delta}(y)=N^{D_{t}^{\delta}}(y)\left(-\frac{d \mu\left(D_{s}^{\delta}\right)}{\underline{\mu}\left(\partial D_{s}^{\delta}\right)}+\frac{1}{2} h^{D_{t}^{\delta}}(y) d t\right) \tag{3.13}
\end{equation*}
$$

Again by constancy in law:

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

So to prove our result we only need to prove that

$$
\begin{equation*}
\int_{0}^{t} \frac{d \mu\left(D_{s}\right)}{\underline{\mu}\left(\partial D_{s}\right)}=-W_{t}+\int_{0}^{t} h^{D_{s}}\left(X_{s}\right) d s+\int_{0}^{t} 2 \sin \left(\theta^{S_{s}}\left(X_{s}\right)\right) d L_{s}^{S_{s}}(X) \tag{3.14}
\end{equation*}
$$

and that

$$
\begin{equation*}
W_{t}=\int_{0}^{t}\left\langle N^{D_{s}}\left(X_{s}\right), d X_{s}\right\rangle \tag{3.15}
\end{equation*}
$$

Let us prove (3.15). In all this paragraph we consider $M$ as isometrically embedded in some Euclidean space. In particular we are allowed to integrate vector 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 G. 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 (G.1) 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 (3.14). It will be a consequence of the convergence as $\delta \rightarrow 0$ 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 $f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)$ :

$$
\begin{align*}
& d\left(f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\right)=\left\langle d_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right), d X_{t}^{\delta}\right\rangle+ \frac{1}{2} \Delta_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right) d t  \tag{3.16}\\
&+\left\langle d_{2} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right), d \partial D_{t}^{\delta}\right\rangle+\frac{1}{2} \nabla_{2} d_{2} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\left(d \partial D_{t}^{\delta}, d \partial D_{t}^{\delta}\right) d t \\
&+\left\langle\nabla_{2} d_{1} f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right), d \partial D_{t}^{\delta} \otimes d X_{t}^{\delta}\right\rangle
\end{align*}
$$

From Proposition 3.4, possibly by extracting a subsequence,

$$
\begin{equation*}
\left(f_{\delta}\left(X_{t}^{\delta}, D_{t}^{\delta}\right)\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(\rho\left(X_{t}, \partial D_{t}\right)\right)_{t \geqslant 0} . \tag{3.17}
\end{equation*}
$$

From (3.5) we get for $i=1,2$,

$$
\begin{aligned}
d_{i} f_{\delta}(x, D)- & d_{i} \rho_{\delta}(x, \partial D) \\
=- & \left.d_{i} \ell_{\varepsilon}(x, D) \int_{T_{x} M} \varphi_{\delta}(|u|)\left(\rho_{\delta}\left(\exp _{x}(u)\right), \partial D\right)-\rho_{\delta}(x, \partial D)\right) d u \\
& \left.+\left(1-\ell_{\varepsilon}(x, D)\right) \int_{T_{x} M} \varphi_{\delta}(|u|)\left(d_{i} \rho_{\delta}\left(\exp _{x}(u)\right), \partial D\right)-d_{i} \rho_{\delta}(x, \partial D)\right) d u .
\end{aligned}
$$

From this we see that $d_{1} f_{\delta}(\cdot, D)$ converges, locally uniformly outside $S(D)$, to $d_{1} \rho(\cdot, \partial D)$ with respect to the distance $d_{0}$ of Appendix G. We obtain, with Lemma G.1, possibly by again extracting a subsequence, that

$$
\begin{equation*}
\left(\int_{0}^{t}\left\langle d_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d X_{s}^{\delta}\right\rangle\right)_{t \geqslant 0} \stackrel{\mathscr{L}}{\longrightarrow}\left(\int_{0}^{t}\left\langle d_{1} \rho\left(X_{s}, \partial D_{s}\right), d X_{s}\right\rangle\right)_{t \geqslant 0} \tag{3.18}
\end{equation*}
$$

More precisely, we have a sequence of martingales converging in law to a martingale $M_{t}$ which is a Brownian motion by [25, Th. 3]. For identifying the limiting martingale, we use the convergence of $\left\langle d_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d X_{s}^{\delta}\right\rangle \otimes d X_{s}^{\delta}$ to $d M_{s} \otimes d X_{s}$ obtained again by [25, Th. 3] (here again we use an isometric embedding of $M$ ). But Lemma G. 1 proves that the limit is equal to $\nabla_{1} \rho\left(X_{s}, \partial D_{s}\right) d s$, yielding (3.18).

Next we prove that

$$
\begin{equation*}
\left(\int_{0}^{t}\left\langle d_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d \partial D_{s}^{\delta}\right\rangle\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(\int_{0}^{t}\left\langle d_{2} \rho\left(X_{s}, \partial D_{s}\right), d \partial D_{s}\right\rangle\right)_{t \geqslant 0} \tag{3.19}
\end{equation*}
$$

The argument is similar except that as we see with (3.8), 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

$$
\begin{align*}
& \left\langle d_{2} f_{\delta}(x, D),\left.N\right|_{\partial D}\right\rangle  \tag{3.20}\\
& \quad=c_{\delta} \int_{T_{x} M} \varphi_{\delta}(|u|)\left\langle-N\left(P\left(\exp _{x}(u)\right), N\left(P\left(\exp _{x}(u)\right)\right\rangle d u=-c_{\delta}\right.\right.
\end{align*}
$$

where $c_{\delta}$ is defined in (3.2). This together with (3.13) suggests to write

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

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

$$
\int_{0}^{t} \tilde{\ell}_{\varepsilon}\left(X_{s}^{\delta}, D_{s}^{\delta}\right)\left\langle d_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right)+c_{\delta} N^{D_{s}^{\delta}}, d \partial D_{s}^{\delta}\right\rangle
$$

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

With this last integral we can proceed as for (3.18), after passing to the limit, and since $\lim _{\delta \rightarrow 0} c_{\delta}=1$, we get (3.20).

Similarly we obtain the two following convergences for the second derivatives:

$$
\begin{align*}
\left(\int_{0}^{t} \nabla_{2} d_{2} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right)\left(d \partial D_{s}^{\delta}, d \partial D_{s}^{\delta}\right)\right)_{t \geqslant 0}  \tag{3.21}\\
\quad \xrightarrow{\mathscr{L}}\left(\int _ { 0 } ^ { t } \nabla _ { 2 } d _ { 2 } \rho ( X _ { s } , \partial D _ { s } ) \left(N\left(P^{\partial D_{s}}\left(X_{s}\right), N\left(P^{\partial D_{s}}\left(X_{s}\right)\right) d s\right)_{t \geqslant 0} \equiv 0\right.\right.
\end{align*}
$$

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),

$$
\begin{align*}
&\left.\left(\int_{0}^{t}\left\langle\nabla_{2} d_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right), d \partial D_{t}^{\delta} \otimes d X_{t}^{\delta}\right\rangle\right)\right)_{t \geqslant 0}  \tag{3.22}\\
& \xrightarrow{\mathscr{L}}\left(\int_{0}^{t}\left\langle\nabla_{2} d_{1} \rho\left(X_{s}, \partial D_{s}\right), d \partial D_{s} \otimes d X_{s}\right\rangle\right)_{t \geqslant 0} \equiv 0
\end{align*}
$$

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 the Itô-Tanaka formula (see Proposition E. 1 in Appendix E using that $\rho(x, \partial D)$ is almost everywhere the minimum of two smooth functions) together with Assumption 3.1 which allows to only consider the regular skeleton, together with Theorem B. 1 which says that the latter has absolutely continuous variation (useful for the term $d L_{t}^{S_{t}}(X)$ ), we have
(3.23) $d\left(\rho\left(X_{t}, \partial D_{t}\right)\right)$

$$
\begin{aligned}
=\left\langle d_{1} \rho\left(X_{t}, \partial D_{t}\right), d X_{t}\right\rangle-\frac{1}{2} h^{D_{t}}\left(X_{t}\right) \mathbb{1}_{D_{t} \backslash S_{t}}\left(X_{t}\right) d t & +\left\langle d_{2} \rho\left(X_{t}, \partial D_{t}\right), d \partial D_{t}\right\rangle \\
+ & +0-\sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X) .
\end{aligned}
$$

Using (3.16), (3.17), (3.18), (3.19), (3.21), (3.22), (3.23) we obtain that

$$
\begin{align*}
& \left(\int_{0}^{t} \Delta_{1} f_{\delta}\left(X_{s}^{\delta}, D_{s}^{\delta}\right) d s\right)_{t \geqslant 0}  \tag{3.24}\\
& \quad \stackrel{\mathscr{L}}{ }\left(\int_{0}^{t}-h^{D_{s}}\left(X_{s}\right) \mathbb{1}_{D_{s} \backslash S_{s}}\left(X_{s}\right) d s-\int_{0}^{t} 2 \sin \left(\theta^{S_{s}}\left(X_{s}\right)\right) d L_{s}^{S_{s}}(X)\right)_{t \geqslant 0}
\end{align*}
$$

It remains to pass in the limit as $\delta$ goes to zero in (3.12), to deduce (3.14).
Remark 3.7. - From (3.23), it can be deduced that

$$
\begin{aligned}
& d\left(\rho\left(X_{t}, \partial D_{t}\right)\right) \\
& \quad=\frac{1}{2}\left(h^{D_{t}}\left(X_{t}\right) \mathbb{1}_{D_{t} \backslash S_{t}}\left(X_{t}\right)-h^{D_{t}}\left(P^{\partial D_{t}}\left(X_{t}\right)\right)\right) d t+\sin \left(\theta^{S_{t}}\left(X_{t}\right)\right) d L_{t}^{S_{t}}(X) .
\end{aligned}
$$

Indeed, (3.15) implies that

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

and due to (3.19), we have

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

where we used (3.12) in conjunction with (3.13).
Taking into account (3.24), we identify the last limit with

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

## 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

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

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

Remark 4.2. - Equation (4.1) can be considered as a limiting case of (2.6). 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 (3.6).

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

Theorem 4.3. - Under Assumption 3.1, there exists an intertwining $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ stopped at $\tau_{\varepsilon}$, satisfying

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

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 $D_{0}:=[-\varepsilon, \varepsilon]$ and next we should let $\varepsilon$ go to $\left.0_{+}\right)$. Due to the invariance by symmetry of (3.6), 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 (3.6) 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 $\left(L_{t}\right)_{t \geqslant 0}$ is the local time of $X$ at 0 . Namely we get that

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

by Tanaka's formula. It is well-known that $\left(R_{t}\right)_{t \geqslant 0}$ is a Bessel process of dimension 3 (see e.g. [19, Chap. 6, Cor.3.8]). In particular the signed distance $\rho_{\partial D_{t}}^{+}$to $\partial D_{t}$ (chosen to be positive inside $D_{t}$ ) is given by

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

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 [18], which is a.s. touching (the upper) boundary repeatedly. Instead we end up with the intertwining dual constructed in [16] 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 (2.6) as $\varepsilon$ goes to zero, on the one hand we have

$$
\mathbb{1}_{\left\{X_{t} \neq R_{t}\right\}} d R_{t}=d X_{t},
$$

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

$$
R_{t}+X_{t}=2 X_{t}-2 \min _{0 \leqslant s \leqslant t} X_{s}
$$

which is equivalent to

$$
R_{t}=X_{t}-2 \min _{0 \leqslant s \leqslant t} X_{s} .
$$

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

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

5.2. Brownian motion and disks in rotationally symmetric manifolds. - This is the simplest 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

$$
d s^{2}=d r^{2}+f^{2}(r) d \Theta^{2}
$$

the metric in polar coordinates. Then the radial Laplacian is

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

We will investigate set-valued processes $\left(D_{t}=B\left(o, R_{t}\right)\right)_{t \geqslant 0}$ 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 $\left(X_{t}\right)_{t \geqslant 0}$ 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

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

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

$$
\mathscr{U}^{f}(d r):=\frac{f(r)}{\int_{0}^{r_{0}} f(s) d s} d r .
$$

The evolution equation (3.6) 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 $\left(R_{t}\right)_{t \geqslant 0}$. Moreover it writes

$$
\begin{align*}
d \rho_{t} & =d \beta_{t}+\frac{1}{2} b\left(\rho_{t}\right) d t  \tag{5.2}\\
d R_{t} & =d \beta_{t}+\left[-\frac{1}{2} b\left(R_{t}\right)+b\left(\rho_{t}\right)\right] d t
\end{align*}
$$

Proposition 5.1. - The system of equations (5.2) has a solution up to explosion time of $\left(R_{t}\right)_{t}$

$$
\tau^{D}:=\inf \left\{t \geqslant 0, R_{t} \notin(0, \infty)\right\}
$$

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

$$
\begin{equation*}
0<\rho_{t}<R_{t} . \tag{5.3}
\end{equation*}
$$

The corresponding set-valued process $\left(D_{t}=B\left(o, R_{t}\right)\right)_{t \geqslant 0}$ is solution to equation (3.6), and in particular, for all $\mathscr{F}^{D}$-stopping time $\tau$,

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

Proof. - We only have to check (5.3). By (5.2),

$$
d\left(R_{t}-\rho_{t}\right)=\frac{1}{2}\left[b\left(\rho_{t}\right)-b\left(R_{t}\right)\right] d t,
$$

which vanishes on $\left\{R_{t}=\rho_{t}\right\}$, and since $b$ is smooth, if $\rho_{0}<R_{0}$, then $\rho_{t}<R_{t}$ for all times.
5.3. Brownian motion and annulus in 2 -dimensional rotationally symmetric manifolds. - 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

$$
d s^{2}=d r^{2}+f^{2}(r) d \theta^{2}
$$

the metric in polar coordinates. Then the radial Laplacian is

$$
\Delta_{r}=\frac{\partial^{2}}{(\partial r)^{2}}+b(r) \frac{\partial}{\partial r} \quad \text { with } \quad b=(\ln f)^{\prime}
$$

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

$$
A\left(r^{-}, r^{+}\right):=\left\{x \in M, r^{-} \leqslant r(x) \leqslant r^{+}\right\} \quad \text { if } \quad r^{-}<r^{+}, \quad A\left(r^{-}, r^{+}\right):=\varnothing
$$

the closed 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

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

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

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

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

$$
\mathscr{U}^{f}\left(\left(r_{0}^{-}, r_{0}^{+}\right)\right)(d r):=\frac{f(r)}{\int_{r_{0}^{-}}^{r_{0}^{+}} f(s) d s} d r
$$

The evolution equation (3.6) for $\left(D_{t}\right)_{t \geqslant 0}$ 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

$$
\begin{align*}
d \rho_{t} & =\operatorname{sign}\left(\rho_{t}-R_{t}^{0}\right) d W_{t}+\frac{1}{2} b\left(\rho_{t}\right) d t, \\
d R_{t}^{+} & =d W_{t}+\left[-\frac{1}{2} b\left(R_{t}^{+}\right)+\operatorname{sign}\left(\rho_{t}-R_{t}^{0}\right) b\left(\rho_{t}\right)\right] d t+2 L_{t}^{R_{t}^{0}}(\rho), \\
d R_{t}^{-} & =-d W_{t}+\left[-\frac{1}{2} b\left(R_{t}^{-}\right)-\operatorname{sign}\left(\rho_{t}-R_{t}^{0}\right) b\left(\rho_{t}\right)\right] d t-2 L_{t}^{R_{t}^{0}}(\rho),  \tag{5.4}\\
R_{t}^{0} & =\frac{1}{2}\left(R_{t}^{-}+R_{t}^{+}\right),
\end{align*}
$$

and these equations imply

$$
d R_{t}^{0}=-\frac{1}{4}\left[b\left(R_{t}^{+}\right)+b\left(R_{t}^{-}\right)\right] d t
$$

Proposition 5.2. - The system of equations (5.4) has a solution up to explosion time

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

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

$$
R_{t}^{-} \leqslant \rho_{t} \leqslant R_{t}^{+}
$$

The corresponding set-valued process $\left(D_{t}=A\left(R_{t}^{-}, R_{t}^{+}\right)\right)_{t \geqslant 0}$ is solution to equation (3.6), and in particular, for all $\mathscr{F}^{D}$-stopping time $\tau$,

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

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 $-\delta / 2<r<\delta / 2$ :

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

and let

$$
\begin{aligned}
\operatorname{sign}_{\delta}: \mathbb{R} & \longrightarrow \mathbb{R} \\
r & \longmapsto-1+2 \int_{-\infty}^{r} \varphi_{\delta}(s) d s
\end{aligned}
$$

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

$$
\begin{aligned}
f^{A\left(r^{-}, r^{+}\right)}(x) & =f\left(x, r^{-}, r^{+}\right) \\
\text {with } \quad g(r) & =g\left(r, r^{-}, r^{+}\right)
\end{aligned}=\int_{r^{-}}^{r}-\operatorname{sign}_{\delta}\left(s-r^{0}\right) d s .
$$

Clearly $f\left(x, r^{-}, r^{+}\right)$is 1-Lipschitz in the first variable. A computation shows that

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

and thus $g$ 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

$$
\langle\nabla f, N\rangle \equiv 1 \quad \text { and } \quad \nabla d f(N, N) \equiv 0
$$

This yields an elementary proof of the properties of Proposition 3.4. We can use Theorem 3.5 to solve equation (5.4) 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 [8, Th. 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 [15] 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 [15] 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. For any integer $n \geqslant 2$, let $G_{n}$ the group of isometries of $\mathbb{R}^{2}$ generated by the rotation of angle $2 \pi / n$ and the symmetry with respect to the horizontal axis. Consider a smooth strictly convex bounded set $\widetilde{D}_{0} \subset M$ with smooth boundary, stable by the action of $G_{n}$. Let us investigate the evolution of $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ solution to (2.7). 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, $\widetilde{D}_{t}$ will stay convex. It will also stay symmetric. All the results of this subsection are proved in [2].
Proposition 5.4. - Assume that its skeleton has the form $\widetilde{S}_{0}=G_{n} \widetilde{H}_{0}, \widetilde{H}_{0}$ being an horizontal interval $\widetilde{H}_{0}=\left[0, \widetilde{x}_{0}\right] \times\{0\}$ for some $\widetilde{x}_{0}>0$ (an example of such a set when $n=2$ is the interior of an ellipse, the skeleton being the interval between the two foci). The skeleton of $\widetilde{D}_{t}$ always takes the form $\widetilde{S}_{t}=G_{n} \widetilde{H}_{t}$ with $\widetilde{H}_{t}=\left[0, \widetilde{x}_{t}\right] \times\{0\}$ an horizontal interval.

The right endpoint $\left(\widetilde{x}_{t}, 0\right)$ in the horizontal axis of the skeleton $\widetilde{S}_{t}$ satisfies

$$
\frac{d \widetilde{x}_{t}}{d t}=\frac{\rho^{2}\left(\left(\widetilde{x}_{t}, 0\right), \widetilde{y}_{t}\right)}{2}\left(h^{\widetilde{D}_{t}}\right)^{\prime \prime}\left(\widetilde{y}_{t}\right),
$$

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

Let us return to the general situation of $G_{n}$-symmetric $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$. The investigation of the lifetime of the solution to (2.7) is not easy. In [2] we prove that the lifetime is the time when $\widetilde{D}_{t}$ meets its skeleton $\widetilde{S}_{t}$. We have no example where $\widetilde{D}_{t}$ meets its skeleton $\widetilde{S}_{t}$ in finite time. The next proposition yields examples where the lifetime is infinite, together with nice properties related to the symmetry group $G_{n}$.

Proposition 5.5
(1) The process $\left(\underline{\mu}^{\partial \widetilde{D}_{t}}\left(\partial \widetilde{D}_{t}\right) / \mu\left(\widetilde{D}_{t}\right)\right)_{0 \leqslant t<\widetilde{\tau}}$ is a supermartingale.
(2) Define the entropy $\widetilde{E n t}_{t}$ as the integral of $\rho_{t} \log \rho_{t}$ with respect to the curvilinear abscissa in $\partial \widetilde{D}_{t}$, $\rho_{t}$ being the curvature of $\partial \widetilde{D}_{t}$. Assume $\widetilde{S}_{0}$ is $G_{n}$-symmetric with $n \geqslant 3$. Then the entropy process $\left(\widetilde{\operatorname{Ent}}_{t}\right)_{0 \leqslant t<\widetilde{\tau}}$ is a supermartingale.
(3) Assume $\widetilde{S}_{0}$ is $G_{n}$-symmetric with $n \geqslant 7$. Then $\widetilde{\tau}=\infty$ a.s. Consequently, when $S_{0}$ is $G_{n}$-symmetric with $n \geqslant 7$, Equation (4.1) for the decoupled $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ provides an intertwining with infinite lifetime. If moreover the skeleton $S_{0}$ of $D_{0}$ has the form $\widetilde{S}_{0}=G_{n} \widetilde{H}_{0}, \widetilde{H}_{0}$ being an horizontal interval $\widetilde{H}_{0}=\left[0, x_{0}\right] \times\{0\}$ for some $x_{0}>0$, Equation (3.6) for the full coupled $\left(X_{t}, D_{t}\right)_{t \geqslant 0}$ provides an intertwining with infinite lifetime.

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

Our goal here is to obtain an extension of Stokes' formula on a domain with a smooth boundary, for functions which degenerate on the skeleton. We take the opportunity to recall this notion, as well as related geometric concepts.

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 (i) the distance to $\partial D$ is not smooth and (ii) there are points around them where the distance to $\partial D$ is smooth with a non vanishing gradient. 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 pair $\left(y_{1}, y_{2}\right)$ 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 with strictly 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 $s, t \in \mathbb{R}$

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

and $\psi(t)=\psi(0, t)$. We will indifferently 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 right-hand side of (A.7) insures that $h$ is non-negative on $\partial D(s)$ when $D(s)$ is convex). This together with (A.6) yields

$$
\left.\frac{d}{d t}\right|_{t=s} \operatorname{det} T \psi(t)(y)=-h(\psi(s)(y)) \operatorname{det} T \psi(s)(y)
$$

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

$$
\operatorname{det} T \psi(t)(y)=\exp \left(\int_{0}^{t}-h(\psi(s)(y)) d s\right)
$$

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

$$
\begin{equation*}
\mu(D)=\int_{0}^{\infty} \underline{\mu}(\partial D(r)) d r \tag{A.8}
\end{equation*}
$$

But for $r \geqslant 0$

$$
\underline{\mu}(\partial D(r))=\int_{\partial D} \operatorname{det} T \psi(r)(y) \underline{\mu}(d y)
$$

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

$$
\underline{\mu}(\partial D(r))=\int_{\partial D} \exp \left(-\int_{0}^{r} h(\psi(s)(y)) d s\right) 1_{\{r<\tau(y)\}} \underline{\mu}(d y),
$$

which yields with (A.8)

$$
\mu(D)=\int_{\partial D}\left(\int_{0}^{\tau(y)} \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right) d r\right) \underline{\mu}(d y)
$$

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

$$
\int_{D} g d \mu=\int_{\partial D}\left(\int_{0}^{\tau(y)} g(\psi(r, y)) \exp \left(-\int_{0}^{r} h(\psi(s, y)) d s\right) d r\right) \underline{\mu}(d y)
$$

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

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

Define the map

$$
\begin{aligned}
\varphi: \partial D & \longrightarrow S^{\prime} \\
y & \longmapsto \psi(\tau(y), y)
\end{aligned}
$$

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 the vector $N\left(\psi\left(\tau\left(y_{i}\right)-, y_{i}\right)\right)$ and the skeleton $S$. 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 $\left.z=\psi\left(\tau\left(y_{i}\right), y_{i}\right)\right) \in S$,

$$
\begin{equation*}
\operatorname{det} T \psi\left(\tau\left(y_{i}\right), y_{i}\right)=\sin \theta\left(\varphi\left(y_{i}\right)\right) \operatorname{det} T \varphi\left(y_{i}\right), \quad i=1,2 \tag{A.9}
\end{equation*}
$$

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 $e_{2}=\nabla \tau(y) /\|\nabla \tau(y)\|$ (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 $\left\langle e_{2}^{\theta}, N(z)\right\rangle<0$ ( $e_{2}^{\theta}$ and $N^{S}(z)$ are not orthogonal, since $\left.\theta(z) \neq \pi / 2\right)$ 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 .


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

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

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

$$
e_{2}^{S}=\cos (\theta(z)) e_{1}^{S}+\sin (\theta(z)) e_{2}^{\theta}
$$

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

$$
\left\langle T \psi\left(e^{\prime}\right), e_{2}^{S}\right\rangle=\sin (\theta(z))\left\langle T \psi\left(e^{\prime}\right), e_{2}^{\theta}\right\rangle .
$$

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

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

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

$$
\left\langle d \tau, e_{2}\right\rangle=\left\langle T \psi\left(e_{2}\right), e_{2}^{\theta}\right\rangle \cot \theta(z) .
$$

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.9).
We arrived at

$$
\int_{D} g h d \mu=\int_{\partial D} g(y) \underline{\mu}(d y)-\int_{\partial D} g(\psi(\tau(y), y)) \operatorname{det} T \psi(\tau(y), y) \underline{\mu}(d y)
$$

$$
+\int_{D}\langle d g, N\rangle d \mu
$$

This yields with (A.9)

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

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:

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

## Appendix B. Moving sets

In this section we describe how to move a domain with smooth boundary by deformation of its boundary. We also 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 $D_{0}$ with smooth boundary $\partial D_{0}$ and $\alpha>0$, define the map $\psi=\psi^{D_{0}}$ by

$$
\begin{aligned}
\psi:(-\alpha, \alpha) \times \partial D_{0} & \longrightarrow M \\
(s, y) & \longmapsto \exp _{y}(s N(y)) .
\end{aligned}
$$

Here $N=N^{D_{0}}$ is the inward normal defined in Section A. We take $\alpha$ sufficiently small so that $\psi$ is a diffeomorphism on its range which we will call $D_{0, \alpha}$. Consider a moving domain $t \mapsto D_{t}$ started at $D_{0}$. We assume that the deformation is sufficiently regular so that for all $t \geqslant 0$, we can write $D_{t}$ as

$$
D_{t}=\left\{\psi\left(\left[f_{t}(y), \tau_{D_{0}}(y)\right], y\right), y \in \partial D_{0}\right\}
$$

with $\tau_{D_{0}}(y)$ defined in (A.1), $\psi\left(\left[f_{t}(y), \tau_{D_{0}}(y)\right], y\right):=\left\{\psi(s, y), s \in\left[f_{t}(y), \tau_{D_{0}}(y)\right]\right\}$, and $t \mapsto f_{t}(y)$ a semimartingale with values in $(-\alpha, \alpha)$, smoothly depending on $y$. In particular, the skeleton $S_{0}^{\prime}$ of $D_{0}$ satisfies $S_{0}^{\prime} \subset D_{t}$. In other words, $D_{t}$ is the union of rays $\psi\left(\left[f_{t}(y), \tau_{D_{0}}(y)\right], y\right)$ orthogonal to $\partial D_{0}$ at $y$ (notice that all $\psi\left(\left[f_{t}(y), \tau_{D_{0}}(y)\right), y\right)$ are disjoint). Alternatively, $D_{t}$ is also the interior of the set $\exp _{\partial D_{0}}(f)$ described in (2.2) with $f_{t}$ instead of $f$. Also, in the special case where the real valued semimartingale $t \mapsto f_{t}(y)=f_{t}$ does not depend on $y$, then we have

$$
\begin{equation*}
D_{t}=D_{0}\left(f_{t}\right) \tag{B.1}
\end{equation*}
$$

where $D_{0}(r)$ is defined in (A.4). 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 the inner skeleton $S_{0}^{\prime}$ or the outer skeleton of $D_{0}$ ), and $t \mapsto f_{t}$ can be allowed to be a semimartingale with singular continuous drift.

When $t \mapsto f_{t}(y)$ depends on $y$ the situation is more complicated and we like to use a more convenient and intrinsic description of the motion of $D_{t}$. More precisely, we will describe it by the motion of its boundary via semimartingales $\left(Y_{t}(y)\right)_{t \geqslant 0}$ indexed by $y \in \partial D_{0}$, satisfying $Y_{0}(y)=y$ and the Itô equation in manifold with respect to the Levi Civita connection $\nabla$

$$
\begin{equation*}
d Y_{t}(y)=d^{\nabla} Y_{t}(y)=N^{D_{t}}\left(Y_{t}(y)\right)\left(H^{D_{t}}\left(Y_{t}(y)\right) d t+d z_{t}\right) \tag{B.2}
\end{equation*}
$$

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

$$
d^{\nabla} Y_{t}(y)=\left(d Y_{t}^{i}(y)+\frac{1}{2} \Gamma_{j, k}^{i}\left(Y_{t}(y)\right) d\left\langle Y_{t}^{j}(y), Y_{t}^{k}(y)\right\rangle\right) D_{i}\left(Y_{t}(y)\right)
$$

where $D_{i}\left(Y_{t}(y)\right)$ is the vector $\partial / \partial y^{i}$ taken at point $Y_{t}(y)$. Since the semimartingale $\left(z_{t}\right)_{t \geqslant 0}$ does not depend on $y$, the Itô equation is equivalent to the Stratonovich one: indeed, using (B.1) the Itô to Stratonovich conversion term is

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

since $N^{D_{t}}\left(Y_{t}(y)\right)$ is the speed at time $a=0$ of the geodesic $a \mapsto \psi^{D_{t}}\left(a, Y_{t}(y)\right)$.
We assume that Equation (B.2) has a strong solution for all times, possibly by stopping it, and that a.s. for all times the map $y^{\prime} \mapsto Y_{t}\left(y^{\prime}\right)$ is a diffeomorphism from $\partial D_{0}$ to $\partial D_{t}$. Since $d Y_{t}(y)$ represents the motion of $\partial D_{t}$, writing $Y_{t}\left(y^{\prime}\right)=y$ and using the diffeomorphism property, equation (B.2) rewrites as

$$
\begin{equation*}
d \partial D_{t}(y):=d Y_{t}\left(y^{\prime}\right)=N^{D_{t}}(y)\left(H^{D_{t}}(y) d t+d z_{t}\right) \tag{B.3}
\end{equation*}
$$

In other words, our equations are driven by two vector fields $\left(H^{D}(y) N^{D}(y)\right)_{y \in \partial D}$ and $\left(N^{D}(y)\right)_{y \in \partial D}$, and the stochastic part is in front of the second one. All the set-valued processes considered in this paper satisfy this assumption.

We can obtain the random functions $f_{t}: \partial D_{0} \rightarrow \mathbb{R}$ from the semimartingales $\left(Y_{t}\left(y^{\prime}\right)\right)_{t \geqslant 0}$ with the following procedure. The orthogonal projection $\pi_{t}: \partial D_{t} \rightarrow \partial D_{0}$ is a diffeomorphism, and by definition of $\psi$, we have

$$
Y_{t}\left(y^{\prime}\right)=\psi\left(f_{t}\left(\pi_{t}\left(Y_{t}\left(y^{\prime}\right)\right), \pi_{t}\left(Y_{t}\left(y^{\prime}\right)\right)\right)\right.
$$

yielding

$$
f_{t}\left(\pi_{t}\left(Y_{t}\left(y^{\prime}\right)\right)=\left(\psi^{-1}\right)_{1}\left(Y_{t}\left(y^{\prime}\right)\right)\right.
$$

with $\left(\psi^{-1}\right)_{1}$ the first coordinate of $\psi^{-1}$. Writing $y=\pi_{t}\left(Y_{t}\left(y^{\prime}\right)\right)$ and using the diffeomorphism properties, we get

$$
f_{t}(y)=\left(\psi^{-1}\right)_{1}\left(\pi_{t}^{-1}(y)\right)
$$

Consequently, the real-valued semimartingale $\left(f_{t}(y)\right)_{t \geqslant 0}$ solves the Stratonovich equation

$$
\circ d f_{t}(y)=T\left(\psi^{-1}\right)_{1}\left(\circ d \pi_{t}^{-1}(y)\right)
$$

which is impossible to work with. This is why we will always consider the formulation (B.3).

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 $J_{i}(0) \in T_{y_{i}} M, J_{1}(1)=J_{2}(1) \in T_{x} M$,

$$
J_{i}(0)=\lambda_{i} N\left(y_{i}\right)+J_{i}^{\perp}(0), \quad J_{i}^{\prime}(0)=\lambda_{i}^{\prime} N\left(y_{i}\right)+\left(J_{i}^{\perp}\right)^{\prime}(0),
$$

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

$$
\begin{equation*}
N_{1}^{S}(x)=\frac{1}{2 \sin \theta(x)}\left(N_{1}(x)-N_{2}(x)\right) . \tag{B.4}
\end{equation*}
$$

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 $N_{1}(x)$.


Figure 2. The points $x, y_{1}, y_{2}$ and the vectors $N_{1}(x), N_{2}(x), N_{1}^{S}(x)$.
We will consider variations of geodesics with same final value:

$$
J_{1}(1)=J_{2}(1)=\lambda N_{1}^{S}(x)+J_{1}^{T}(1)
$$

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

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

and

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

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

$$
\left\langle J_{1}(1), N_{1}(x)\right\rangle-\left\langle J_{1}(0), N\left(y_{1}\right)\right\rangle=\left\langle J_{2}(1), N_{2}(x)\right\rangle-\left\langle J_{2}(0), N\left(y_{2}\right)\right\rangle
$$

or

$$
\lambda \sin \theta(x)+\left\langle J_{1}^{T}(1), N_{1}(x)\right\rangle-\lambda_{1}=-\lambda \sin \theta(x)+\left\langle J_{1}^{T}(1), N_{2}(x)\right\rangle-\lambda_{2},
$$

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

$$
\begin{equation*}
\left\langle J_{1}(1), N_{1}^{S}(x)\right\rangle N_{1}^{S}(x)=\frac{\lambda_{1}-\lambda_{2}}{2 \sin \theta(x)} N_{1}^{S}(x) \tag{B.5}
\end{equation*}
$$

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

$$
J_{1}(0)=\lambda_{1} N\left(y_{1}\right), \text { i.e., } \quad J_{1}^{\perp}(0)=0 .
$$

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

$$
J_{1}^{\perp}(1)=J\left(1,0,\left(J_{1}^{\perp}\right)^{\prime}(0)\right),
$$

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

$$
\begin{aligned}
& J_{1}(1)=J_{1}^{T}(1)+\left\langle J_{1}(1), N_{1}^{S}(x)\right\rangle N_{1}^{S}(x), \\
& J_{1}(1)=J_{1}^{\perp}(1)+\left\langle J_{1}(1), N_{1}(x)\right\rangle N_{1}(x),
\end{aligned}
$$

we get

$$
\begin{equation*}
J_{1}^{T}(1)=J_{1}^{\perp}(1)+\left\langle J_{1}(1), N_{1}(x)\right\rangle N_{1}(x)-\left\langle J_{1}(1), N_{1}^{S}(x)\right\rangle N_{1}^{S}(x) \tag{B.6}
\end{equation*}
$$

On the other hand we have

$$
\begin{aligned}
& \left\langle J_{1}(1), N_{2}(x)\right\rangle=\left\langle J_{1}^{\perp}(1), N_{2}(x)\right\rangle+\left\langle J_{1}(1), N_{1}(x)\right\rangle\left\langle N_{1}(x), N_{2}(x)\right\rangle, \\
& \left\langle J_{1}(1), N_{2}(x)\right\rangle=\left\langle J_{1}(1), N_{1}(x)\right\rangle-\left(\lambda_{1}-\lambda_{2}\right),
\end{aligned}
$$

where the second equation is a direct consequence of (B.5). Subtracting the second equation to the first one yields

$$
(1-\cos (2 \theta(x)))\left\langle J_{1}(1), N_{1}(x)\right\rangle=\left\langle J_{1}^{\perp}(1), N_{2}(x)\right\rangle+\lambda_{1}-\lambda_{2}
$$

Replacing $\left\langle J_{1}(1), N_{1}(x)\right\rangle$ in (B.6) and after simplification, using (B.4) and (B.5), we finally obtain the horizontal displacement

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

We are now in position to write the motion of the skeleton $S_{t}$ when the motion of the boundary is given by (B.3). For $x \in S_{t}$ with corresponding points $y_{1}$ and $y_{2}$ in $\partial D_{t}$,

$$
\begin{equation*}
d S_{t}^{\perp}(x)=\frac{1}{2 \sin \theta^{S_{t}}(x)}\left(H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)\right) N_{1}^{S_{t}}(x) d t \tag{B.7}
\end{equation*}
$$

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.3) we deduce the rotation

$$
\left(J_{1}^{\perp}\right)^{\prime}(0) d t=\rho_{S}\left(y_{1}\right) \nabla_{t} N^{D_{t}}\left(y_{1}\right)=-\rho_{S}\left(y_{1}\right) \nabla H^{D_{t}}\left(y_{1}\right) d t
$$

(in the right-hand side 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 $J_{1}^{T}(1) d t$

$$
\begin{align*}
& J_{1}^{T}(1) d t=J_{1}^{\perp}(1) d t+\frac{1}{4 \sin ^{2} \theta^{S_{t}}(x)}\left(2\left\langle J_{1}^{\perp}(1), N_{2}^{D_{t}}(x)\right\rangle N_{1}^{D_{t}}(x)\right.  \tag{B.8}\\
&\left.+\left(H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)\right)\left(N_{1}^{D_{t}}(x)+N_{2}^{D_{t}}(x)\right)\right) d t
\end{align*}
$$

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 process $z_{t}$ 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.3)

$$
\begin{equation*}
d \partial D_{t}(y)=N^{D_{t}}(y)\left(H^{D_{t}}(y) d t+d z_{t}\right) \tag{B.9}
\end{equation*}
$$

the regular skeleton $S_{t}$ has the normal evolution (B.7)

$$
\begin{equation*}
d S_{t}^{\perp}(x)=\frac{H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)}{4 \sin ^{2} \theta^{S_{t}}(x)}\left(N_{1}^{D_{t}}(x)-N_{2}^{D_{t}}(x)\right) d t \tag{B.10}
\end{equation*}
$$

and the tangential evolution (B.8) which can be rewritten as

$$
\begin{align*}
d S_{t}^{T}(x) & =p_{S}\left(J_{1}^{\perp}(1)\right) d t  \tag{B.11}\\
& +\left(-\frac{\left\langle J_{1}^{\perp}(1), N_{1}^{S}(x)\right\rangle}{2 \sin \theta^{S_{t}}(x)}+\frac{H^{D_{t}}\left(y_{1}\right)-H^{D_{t}}\left(y_{2}\right)}{4 \sin ^{2} \theta^{S_{t}}(x)}\right)\left(N_{1}^{D_{t}}(x)+N_{2}^{D_{t}}(x)\right) d t,
\end{align*}
$$

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.10) is symmetric in $y_{1}$ and $y_{2}$, formula (B.11) 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.9), 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 (2.7)

In this section we adapt the results of [8] to our notations. Let the stochastic mean curvature flow be a solution of:

$$
\begin{equation*}
\forall t \in[0, \tau), \forall y \in C_{t}, \quad d \partial D_{t}(y)=\left(d W_{t}+\frac{1}{2} h^{D_{t}}(y) d t\right) N^{D_{t}}(y) \tag{C.1}
\end{equation*}
$$

where $C_{t}:=\partial D_{t}$, starting at $D_{0}$. Notice that contrarily to [8] we don't have a term $\sqrt{2}$ in front of the Brownian motion, this explains the fact that we have put a normalization factor $1 / 2$ in front of the mean curvature term.

Let $\partial G_{t}$ be a solution of

$$
\left\{\begin{array}{l}
G_{0}=D_{0},  \tag{C.2}\\
\partial_{t} x=\alpha_{\partial G_{t},-W_{t}}(x) N^{G_{t}}(x), \quad \forall t \in[0, \widetilde{\varepsilon}), \forall x \in \partial G_{t},
\end{array}\right.
$$

for some $\widetilde{\varepsilon}>0$ small enough, where $\alpha$ is defined by

$$
\forall r>0, \forall D \in \mathcal{D}_{r}, \forall x \in C, \quad \alpha_{C, r}(x):=\frac{1}{2} h^{\Psi(C, r)}\left(\psi_{C, r}(x)\right)
$$

and $\Psi(C, r)$ is the normal (exterior) flow starting at $C$ at time $r$ (cf. [8, Chap. 3 \& 4] for the notations).

Similarly to the proof of $[8, \mathrm{Th} .17]$, 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 $\mathbb{F}_{f}(D)=\int_{D} f d \mu, \nu=-N$ is the exterior normal

$$
\widetilde{\mathcal{L}} \mathbb{F}_{f}(D):=\frac{1}{2} \int_{\partial D}\langle\nabla f, \nu\rangle d \underline{\mu}=\mathbb{F}_{\frac{1}{2} \Delta f}(D)
$$

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 [8, §4.1], we show that the solution of (C.2) have a regularity $C^{1+\alpha / 2,2+\alpha}$, 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:

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

where in the first equality we use the Itô formula, the fact that $t \mapsto \partial G_{t}$ is of class $C^{1+\alpha / 2},\left(d^{2} / d^{2} r\right) \Psi(x, r)=0$, and in the second equality we used [8, Lem. 13], i.e., $\partial D_{t}$ is a solution in the Itô form:

$$
\left\{\begin{array}{l}
d \partial D_{t}(x)=\left(d W_{t}+\frac{1}{2} h^{\partial D_{t}}(x) d t\right) N^{\partial D_{t}}(x)  \tag{C.3}\\
x \in \partial D_{t}
\end{array}\right.
$$

Proposition C.2. - Conversely, if $\partial D_{t}$ is a solution of (C.3) then $\partial G_{t}=\Psi\left(\partial D_{t}, W_{t}\right)$ is a solution of (C.2).

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

where we use that in this case, the Stratonovich differential is equal to the Itô's one (cf. Appendix B), i.e., $\circ d \partial D_{t}(x)=d \partial D_{t}$, and $\left(d^{2} / d^{2} r\right) \Psi(x, r)=0$. So $\partial G_{t}$ is a solution of (C.2).

By the uniqueness of the solution of (C.2) (cf. [8, Th. 22]) and the fact that it is adapted to the filtration of $B$ we deduce that the solution of (C.3) 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 (2.7) 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 [22] or Ethier and Kurtz [10]. 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

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

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

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

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 $\mathcal{R}$ of "regular" functions that will be important for our purposes is that defined by

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

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

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

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

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

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

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

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

$$
L: \mathcal{D}(L) \longrightarrow \mathcal{R}
$$

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

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

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 $g \in \mathcal{R}$,

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

(we came back to the usual convention that $t>0$ in the above limit) and as a byproduct, 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

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

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 that (D.2) holds.

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 left-hand side is equal to

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

where $y=X_{s}$. By Fubini's lemma, the previous right-hand side can be written as

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

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 Borel 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 $\widetilde{X}:=\left(\widetilde{X}_{t}\right)_{t \geqslant 0}$ 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 $\widetilde{f}$ coinciding with $f$ on $V \backslash A$. The domain $\mathcal{D}(\widetilde{L})$ is the set of $\widetilde{f} \in \widetilde{\mathcal{R}}$ such that there exists $f \in \mathcal{D}(L)$ with $\widetilde{f}$ 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 [21, Th. 7.29], the authors are working with Euclidean spaces, but the result can be extended to metric spaces).

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

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

Such relations are not so obvious if we had chosen to work in a Banach setting (cf. e.g. the book of Yosida [24]), 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 $\mathcal{C}_{b}(V)$.

## Appendix E. 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 E.1. - Let $X_{t}$ a Brownian motion in $M$. We have the following ItôTanaka 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 (3.7).

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

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

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

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

After changing the role of $f$ and $g$ we get

$$
\begin{align*}
d \rho_{\partial D}^{+}\left(X_{t}\right)=\frac{1}{2} \mathbb{1}_{X_{t} \notin S} \Delta \rho_{\partial D}^{+}\left(X_{t}\right) d t &  \tag{E.1}\\
& \quad+\mathbb{1}_{X_{t} \notin S}\left\langle\nabla \rho_{\partial D}^{+}\left(X_{t}\right), d X_{t}\right\rangle-\frac{1}{2} d L_{t}^{0}((f-g)(X .)),
\end{align*}
$$

where

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

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 $(d / d t) \gamma(t)=-\nabla(f-g)(\gamma(t)) /\|\nabla(f-g)(\gamma(t))\|^{2}$ that starts at $y \in U$, we get

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

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

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

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 (E.1), this yield the proposition.

## Appendix F. Uniqueness in law of $\widetilde{\mathcal{L}}$ 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 champ 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 $D \in \mathcal{D}^{2+\alpha}$,

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

Note that $\widehat{\mathscr{L}}$ has the same carré du champ as the carré du champ associated to $\widetilde{\mathscr{L}}$. From now the generator $\widehat{\mathscr{L}}$ is defined as in (2.10).

Proposition F.1. - The martingale problem associated $\widehat{\mathscr{L}}$ is well-posed.
Proof. - We have already shown the existence result in [8], 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 $\frac{1}{2} \Delta k_{\lambda}(x, y)=$
$\left(\lambda^{2} / 2\right) k_{\lambda}(x, y)$. 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 $\left(\widehat{D}_{t}\right)_{t \geqslant 0}$ 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

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

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

$$
\begin{equation*}
\mathbb{E}\left[f_{\lambda}\left(X_{t}, D_{0}\right)\right]=\mathbb{E}\left[f_{\lambda}\left(X_{0}, \widehat{D}_{t}\right)\right] \tag{F.1}
\end{equation*}
$$

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 $\left(\widetilde{D}_{t}\right)_{t \geqslant 0}$ that starts at $D_{0}$ :

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

and so

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

In order to apply [10, Th. 4.2], 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 $n \in \mathbb{N}$

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

The above computations could be done also for $\widetilde{k}_{\lambda_{1}, \lambda_{2}}=e^{\lambda_{1} x+\lambda_{2} y}$, since $\frac{1}{2} \Delta \widetilde{k}_{\lambda_{1}, \lambda_{2}}=$ $\left(\lambda_{1}^{2}+\lambda_{2}^{2}\right) / 2 \widetilde{k}_{\lambda_{1}, \lambda_{2}}$, 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 (F.1) 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 F.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 $k \in C^{\infty}(M)$

$$
\begin{equation*}
F_{k}\left(D_{t}\right)=F_{k}\left(D_{0}\right)+\int_{0}^{t} \mathscr{L}\left[F_{k}\right]\left(D_{s}\right) d s+\int_{0}^{t} \sqrt{\Gamma_{\mathscr{L}}\left[F_{k}, F_{k}\right]}\left(D_{s}\right) d B_{s} \tag{F.2}
\end{equation*}
$$

where $\sqrt{\Gamma_{\mathscr{L}}\left[F_{k}, F_{k}\right]}(D):=\int_{\partial D} k d \sigma$, this is actually [8, Prop. 53]. Note that this procedure of enlargement ([19, Chap. V, Th. 1.7]) could be done by gluing the same independent Brownian motion for each $\left(\Omega, \mathcal{F}^{D}, \mathbb{Q}\right)$. We denote by $\left(\widetilde{\Omega}, \widetilde{\mathcal{F}}^{D}, \widetilde{\mathbb{Q}}\right)$ 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 (F.2) becomes in a differential form

$$
d F_{k}\left(D_{t}\right)-\widehat{\mathscr{L}}\left[F_{k}\right]\left(D_{t}\right) d t=\left(\int_{\partial D} k d \sigma\right)\left(d B_{t}+\frac{\underline{\mu}^{\partial D_{t}}\left(\partial D_{t}\right)}{\mu\left(D_{t}\right)} d t\right)
$$

Let

$$
\begin{aligned}
M_{t} & =\exp \left(-\int_{0}^{t}\left\langle\frac{\underline{\mu}^{\partial D_{s}}\left(\partial D_{s}\right)}{\mu\left(D_{s}\right)}, d B_{s}\right\rangle-\frac{1}{2} \int_{0}^{t}\left(\frac{\underline{\mu}^{\partial D_{s}}\left(\partial D_{s}\right)}{\mu\left(D_{s}\right)}\right)^{2} d s\right), \\
\mathbb{P}_{\mid \mathfrak{F}_{t}} & =M_{t} \widetilde{\mathbb{Q}}_{\mid \mathfrak{F}_{t}} .
\end{aligned}
$$

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

## Appendix G. Convergence in law: a key lemma

This Appendix is devoted to the adaptation to some domain-valued sequences of processes, of [25, Lem.4], which states stability of some time integrals under convergence in law.

Lemma G.1. - Let $\widetilde{\mathscr{F}}:=\widetilde{\mathscr{F}}^{\alpha}, \varepsilon$. We endow the $\operatorname{set} \mathscr{C}([0, \infty), M \times \widetilde{\mathscr{F}})$ of continuous paths with the two dissimilarity measures $d_{\beta}, \beta \in\{0, \alpha\}$, defined as:

$$
d_{\beta}\left(\left(x^{1}, D^{1}\right),\left(x^{2}, D^{2}\right)\right)=\sup _{t \geqslant 0} \rho\left(x^{1}(t), x^{2}(t)\right)+\sup _{t \geqslant 0} d_{\beta, \widetilde{\mathscr{F}}}\left(D^{1}(t), D^{2}(t)\right)
$$

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

$$
d_{\beta, \widetilde{\mathscr{F}}}\left(D, D^{\prime}\right)= \begin{cases}d_{\beta, D}\left(D, D^{\prime}\right) \wedge d_{\beta, D^{\prime}}\left(D^{\prime}, D\right) \wedge \varepsilon & \text { if } H\left(D, D^{\prime}\right)<\varepsilon \\ \varepsilon & \text { otherwise }\end{cases}
$$

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 (2.3).

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 (3.10) converging in law to the limit defined in (3.11) 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 $M \times \widetilde{\mathscr{F}}$ with values in some Euclidean space, and $U$ an open set in $M \times \widetilde{\mathscr{F}}$ 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 G.2. - In the applications we will always take

$$
\begin{equation*}
U=\{(x, D) \in M \times \widetilde{\mathscr{F}}, x \in D \backslash S(D)\} \tag{G.1}
\end{equation*}
$$

which is easily seen to be $d_{0}$-open thanks to Assumption 3.1 on $\widetilde{\mathscr{F}}$.
Proof. - We will follow the proof of [25, Lem. 4], but with several differences due to infinite dimensional spaces. Set for $n \in \mathbb{N}, t \geqslant 0$,

$$
\begin{equation*}
A_{t}^{n}:=\int_{0}^{t} f_{n}\left(X_{s}^{n}, D_{s}^{n}\right) d s, \quad A_{t}:=\int_{0}^{t} f\left(X_{s}, D_{s}\right) d s \tag{G.2}
\end{equation*}
$$

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

$$
\left(X_{t}^{n}, D_{t}^{n}, A_{t}^{n}\right)_{t \geqslant 0} \xrightarrow{\mathscr{L}}\left(X_{t}, D_{t}, a_{t}\right)_{t \geqslant 0}
$$

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

$$
\left(X_{t}^{n}, D_{t}^{n}, A_{t}^{n}, X_{t}, D_{t}, a_{t}\right)_{t \geqslant 0}
$$

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

$$
\begin{equation*}
\left(Z_{t}^{n}\right)_{t \geqslant 0}:=\left(X_{t}^{n}, D_{t}^{n}, A_{t}^{n}\right)_{t \geqslant 0} \xrightarrow{\text { a.s. }}\left(X_{t}, D_{t}, a_{t}\right)_{t \geqslant 0}=:\left(Z_{t}\right)_{t \geqslant 0} . \tag{G.3}
\end{equation*}
$$

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

$$
S:=\left\{\left(X_{s}(\omega), D_{s}(\omega)\right), s \in\left[t-\varepsilon^{\prime}, t+\varepsilon^{\prime}\right]\right\}
$$

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

$$
V=\left\{(x, D) \in M \times \widetilde{\mathscr{F}}, d_{\alpha}((x, D), S) \leqslant \varepsilon^{\prime \prime}\right\}
$$

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 Arzelà-Ascoli 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

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

Both terms in the right converge to 0 , the first one by (a) and the second one by (b). So we have by (G.3) and the above calculation

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

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 $t \geqslant 0$,

$$
\frac{d}{d t} a_{t}(\omega)=f\left(X_{t}(\omega), D_{t}(\omega)\right) \quad \text { a.s. }
$$

This implies that $\omega$ a.s.

$$
\begin{equation*}
\frac{d}{d t} a_{t}(\omega)=f\left(X_{t}(\omega), D_{t}(\omega)\right) \quad \text { for a.e. } t \tag{G.4}
\end{equation*}
$$

On the other hand we know by [14, Th. 10] that $\left(a_{t}\right)_{t \geqslant 0}$ is absolutely continuous:

$$
\begin{equation*}
a_{t}(\omega)=\int_{0}^{t} \ell_{s}(\omega) d s \tag{G.5}
\end{equation*}
$$

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

$$
\begin{equation*}
\lim _{\varepsilon \searrow 0} \frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon}\left|\ell_{s}(\omega)-\ell_{t}(\omega)\right| d s=0 \tag{G.6}
\end{equation*}
$$

Equalities (G.4) and (G.5) imply that $\omega$ a.s.

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

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 (G.6) implies that $\omega$ a.s. for a.e. $t \geqslant 0$

$$
\begin{equation*}
\lim _{\varepsilon \searrow 0} \frac{1}{2 \varepsilon} \int_{t-\varepsilon}^{t+\varepsilon} \ell_{s}(\omega) d s=\ell_{t}(\omega) . \tag{G.7}
\end{equation*}
$$

Consequently, using (G.4) and (G.7), we get $\omega$ a.s. for a.e. $t \geqslant 0$

$$
\ell_{t}(\omega)=f\left(X_{t}(\omega), D_{t}(\omega)\right)
$$

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

$$
a_{t}(\omega)=A_{t}(\omega)=\int_{0}^{t} f\left(X_{s}(\omega), D_{s}(\omega)\right) d s
$$

This together with (G.2) proves the lemma.

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