open science

# Stochastic processes associated to multidimensional parabolic transmission problems in divergence form 

Pierre Etoré, Miguel Martinez

## - To cite this version:

Pierre Etoré, Miguel Martinez. Stochastic processes associated to multidimensional parabolic transmission problems in divergence form. ALEA: Latin American Journal of Probability and Mathematical Statistics, 2023, 20, pp.291-311. 10.30757/ALEA.v20-11 . hal-02531467v3

## HAL Id: hal-02531467 <br> https://hal.science/hal-02531467v3

Submitted on 20 Feb 2023

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire HAL, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

# Stochastic processes associated to multidimensional parabolic transmission problems in divergence form 

Pierre Etoré and Miguel Martinez

Univ. Grenoble Alpes, CNRS, Inria, Grenoble INP, LJK 700 Avenue Centrale, 38401 St Martin D'Hères, France, *corresponding author
E-mail address: pierre.etore@univ-grenoble-alpes.fr
URL: https://membres-ljk.imag.fr/Pierre.Etore/
Univ. Gustave Eiffel, LAMA, 5, boulevard Descartes, 77454 Marne-la-Vallee cedex 2, France
E-mail address: miguel.martinez@u-pem.fr
URL: https://perso.math.u-pem.fr/martinez.miguel/


#### Abstract

In this note we define and study the stochastic process $X$ in link with a parabolic transmission operator $(A, \mathcal{D}(A))$ in divergence form. The transmission operator involves a diffraction condition along a transmission boundary. To that aim we gather and clarify some results coming from the theory of Dirichlet forms as exposed in Fukushima et al. (1994) and Stroock (1988) for general divergence form operators. We show that $X$ is a semimartingale and that it is solution of a stochastic differential equation involving partial reflections in the co-normal directions along the transmission boundary.


## 1. Introduction

In this note we aim at tying - with all the necessary rigor - various theoretical results that come from different approaches concerning the probabilistic study of divergence form operators. We also provide a probabilistic representation of the underlying process when the divergence operator is a transmission operator involving a transmission condition across some smooth interface: in this particular case, we show that the process is solution of a stochastic differential equation involving partial reflections in the co-normal directions along the transmission boundary. All of these results are natural but we could not find them in the existing literature and would like to record them in print with a proof as they ground the foundations for the study of probabilistic numerical methods for transmission problems (see e.g. Bossy et al. (2010); Étoré and Martinez (2021)).

Our starting point is $a: \mathbb{R}^{d} \rightarrow \mathcal{M}_{d}(\mathbb{R})$ some measurable symmetric matrix valued coefficient satisfying the following ellipticity and boundedness condition $(\mathbf{E}-\mathbf{B})$ :

Assumption 1.1. (E-B): There exists $\lambda, \Lambda \in(0, \infty)$ such that

$$
\begin{equation*}
\forall x \in \mathbb{R}^{d}, \quad \forall \xi \in \mathbb{R}^{d}, \quad \lambda|\xi|^{2} \leq \xi^{*} a(x) \xi \leq \Lambda|\xi|^{2} . \tag{1.1}
\end{equation*}
$$

[^0]Let us associate to the coefficient $a$ the following unbounded operator $A: \mathcal{D}(A) \subset L^{2}\left(\mathbb{R}^{d}\right) \rightarrow$ $L^{2}\left(\mathbb{R}^{d}\right)$ defined by

$$
\mathcal{D}(A)=\left\{u \in H^{1}\left(\mathbb{R}^{d}\right) \text { with } \sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} u\right) \in L^{2}\left(\mathbb{R}^{d}\right)\right\}
$$

and

$$
\forall u \in \mathcal{D}(A), \quad A u=\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} u\right)
$$

There exists a closed symmetric Dirichlet form $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ and its corresponding semigroup $\left(T_{t}\right)$ on $L^{2}\left(\mathbb{R}^{d}\right)$ that are naturally in link with $(A, \mathcal{D}(A))$.

We define rigorously these objects and study their relations in Section 2. Using the spectral resolution of the identity associated to $(A, \mathcal{D}(A))$, we study the regularity in the 'time variable $t$ ' of $\mathcal{E}\left(T_{t} f, g\right), f \in L^{2}\left(\mathbb{R}^{d}\right), g \in \mathcal{D}[\mathcal{E}]$ (Subsection 2.1). This permits to establish rigorously in Subsection 2.2 the connection with the results in Stroock (1988) that are exposed by D.W. Stroock in the $C_{b}\left(\mathbb{R}^{d}\right)$ setting (Feller semigroup) and to assert the validity of Aronson's estimates for the transition function of ( $T_{t}$ ) (see Aronson (1967), Stroock (1988), Bass (1998)).

We then aim at providing tractable (from a numerical perspective) stochastic representations for the Hunt process $X$ associated to $(\mathcal{E}, \mathcal{D}, \mathcal{E})$.

Of course, we are in the ideal setting to apply the stochastic calculus for symmetric Dirichlet forms, and we give the Fukushima decomposition of $X$ via the Revuz correspondence for additive functionals that is presented in Subsection 3.1. In this general setting a representation of $X$ may be also provided by the so-called Lyons-Zheng decomposition involving reversed-time martingale increments: we give a brief insight of the ideas behind this theoretical decomposition in Subsection 3.2. Please note that none of the results presented in Section 3 are new (see Fukushima et al. (1994), Lyons and Zheng (1990)) and we have tried our best to present the ideas in a coherent and assimilable way for a reader that might not be familiar with the subject.

Then, gradually moving from broad issues to more specific ones, we focus in Section 4 on the particular case where $(A, \mathcal{D}(A))$ is a transmission operator across some transmission boundary: we present a Skorokhod representation of the Hunt process $X$ associated to ( $\mathcal{E}, \mathcal{D}[\mathcal{E}])$ in this case. This result is new and constitutes the main contribution of this note (see Subsection 4.1). Finally, we give a special attention to the particular case of a diagonal coefficient matrix $a$ that remains constant on each side of the transmission boundary and compare our result to the one obtained in the pionneering paper Bossy et al. (2010) (see Subsection 4.2). We show that, when reduced to this very specific context, our description essentially matches the stochastic differential equation considered in Bossy et al. (2010).

## 2. Dirichlet form and Markovian semigroup associated to general elliptic divergence form operators

2.1. Definitions and first properties. To the coefficient matrix $a$, we may associate a closed symmetric Dirichlet form $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ defined on $L^{2}\left(\mathbb{R}^{d}\right)$ by

$$
\left\{\begin{array}{l}
\mathcal{D}[\mathcal{E}]=H^{1}\left(\mathbb{R}^{d}\right), \\
\mathcal{E}(u, v)=\sum_{i, j=1}^{d} \int_{\mathbb{R}^{d}} a_{i j} D_{j} u D_{i} v, \quad u, v \in \mathcal{D}[\mathcal{E}]
\end{array}\right.
$$

(see Fukushima et al. (1994), p111). This closed symmetric Dirichlet form is the starting point of our construction.

On the underlying Hilbert space $L^{2}\left(\mathbb{R}^{d}\right)$, we denote within this subsection by $(A, \mathcal{D}(A))$ the (unique) self-adjoint operator associated to $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ and characterized by

$$
\left\{\begin{array}{l}
\mathcal{D}(A) \subset \mathcal{D}[\mathcal{E}] \\
\mathcal{E}(u, v)=-\langle A u, v\rangle_{L^{2}\left(\mathbb{R}^{d}\right)}, \quad u \in \mathcal{D}(A), v \in \mathcal{D}[\mathcal{E}]
\end{array}\right.
$$

(Fukushima et al. (1994), Theorem 1.3.1 and Corollary 1.3.1 p.21).
We aim at identifying this operator - as expected it will turn out that $(A, \mathcal{D}(A))$ is nothing else than the operator defined in the Introduction, therefore the common notation.

By the very definition of $(A, \mathcal{D}(A))$, we have for any $f \in \mathcal{D}(A)$ and any $g \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$

$$
-\langle A f, g\rangle_{L^{2}\left(\mathbb{R}^{d}\right)}=\mathcal{E}(f, g)=\sum_{i, j=1}^{d} \int_{\mathbb{R}^{d}} a_{i j} D_{j} f D_{i} g=-\left\langle\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right), g\right\rangle_{H^{-1}\left(\mathbb{R}^{d}\right), H^{1}\left(\mathbb{R}^{d}\right)}
$$

where $\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right)$ is understood in the distributional sense as an element of $H^{-1}\left(\mathbb{R}^{d}\right)$. But as $A f \in L^{2}\left(\mathbb{R}^{d}\right)$ by the definition of $\mathcal{D}(A)$ the above equality shows that $\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right) \in L^{2}\left(\mathbb{R}^{d}\right)$ (for any $f \in \mathcal{D}(A)$ ).

Thus, it is proved that $\mathcal{D}(A) \subseteq\left\{f \in H^{1}\left(\mathbb{R}^{d}\right)\right.$ with $\left.\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right) \in L^{2}\left(\mathbb{R}^{d}\right)\right\}$.
In turn (by the density of $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ in $L^{2}\left(\mathbb{R}^{d}\right)$ ) the equality permits to identify for any $f \in \mathcal{D}(A)$,

$$
A f=\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right)
$$

Let us now prove the reverse inclusion $\left\{f \in H^{1}\left(\mathbb{R}^{d}\right)\right.$ with $\left.\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right) \in L^{2}\left(\mathbb{R}^{d}\right)\right\} \subseteq \mathcal{D}(A)$.
Let $f \in\left\{f \in H^{1}\left(\mathbb{R}^{d}\right)\right.$ with $\left.\sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right) \in L^{2}\left(\mathbb{R}^{d}\right)\right\}$. By the symmetry of the coefficient matrix $a$ and integration by parts, it is not hard to prove that for any $v \in \mathcal{D}(A)$,

$$
\langle A v, f\rangle_{L^{2}\left(\mathbb{R}^{d}\right)}=-\mathcal{E}(v, f)=-\sum_{j, i=1}^{d} \int_{\mathbb{R}^{d}} a_{j i} D_{i} f D_{j} v=\left\langle\sum_{j, i} D_{j}\left(a_{j i} D_{i} f\right), v\right\rangle_{L^{2}\left(\mathbb{R}^{d}\right)}
$$

and in particular $f \in \mathcal{D}\left(A^{*}\right) \stackrel{\text { def }}{=}\left\{g \in L^{2}\left(\mathbb{R}^{d}\right) \mid \exists h_{g} \in L^{2}\left(\mathbb{R}^{d}\right)\right.$ s.t. $\left.\langle A v, g\rangle=\left\langle v, h_{g}\right\rangle, \forall v \in \mathcal{D}(A)\right\}$ (see Pazy (1983)). So that we get the reverse inclusion

$$
\left\{f \in H^{1}\left(\mathbb{R}^{d}\right) \text { with } \sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right) \in L^{2}\left(\mathbb{R}^{d}\right)\right\} \subseteq \mathcal{D}\left(A^{*}\right)=\mathcal{D}(A)
$$

where the equality comes from the fact that $(A, \mathcal{D}(A))$ is self-adjoint. Finally, we have proved

$$
\begin{equation*}
\mathcal{D}(A)=\left\{f \in H^{1}\left(\mathbb{R}^{d}\right) \text { with } \sum_{i, j=1}^{d} D_{i}\left(a_{i j} D_{j} f\right) \in L^{2}\left(\mathbb{R}^{d}\right)\right\} \tag{2.1}
\end{equation*}
$$

and $(A, \mathcal{D}(A))$ is fully identified as being the same operator of the Introduction 1.
Note that since $a$ is only assumed to be measurable, $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ - which is a core for the Dirichlet form $(\mathcal{E}, \mathcal{D}(\mathcal{E}))$ - is not even a subset of $\mathcal{D}(A)$.

We now turn to the study of the spectral resolution and the semigroup associated to $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ and $(A, \mathcal{D}(A))$. For the sake of conciseness we denote $(\cdot, \cdot)=\langle\cdot, \cdot\rangle_{L^{2}\left(\mathbb{R}^{d}\right)}$ and $\|\cdot\|=\|\cdot\|_{L^{2}\left(\mathbb{R}^{d}\right)}$ till the end of the section.

Since $(-A, \mathcal{D}(A))$ is a self-adjoint operator on the Hilbert space $L^{2}\left(\mathbb{R}^{d}\right)$ that is non-negative definite, it admits a spectral resolution of the identity $\left\{E_{\gamma}: \gamma \in[0, \infty)\right\}$. For any $\gamma \geq 0$ the operator $E_{\gamma}: L^{2}\left(\mathbb{R}^{d}\right) \rightarrow L^{2}\left(\mathbb{R}^{d}\right)$ is a self-adjoint projection operator with $\left(E_{\gamma} f, f\right) \geq 0, f \in L^{2}\left(\mathbb{R}^{d}\right)$,
and the $E_{\gamma}$ 's form a spectral family with in particular $E_{\mu} E_{\gamma}=E_{\mu \wedge \gamma}$, (see Fukushima et al. (1994) p 18 for a list of properties). The link with $(-A, \mathcal{D}(A))$ is through

$$
(-A f, g)=\int_{[0, \infty)} \gamma d\left(E_{\gamma} f, g\right) \quad \forall f \in \mathcal{D}(A), g \in L^{2}\left(\mathbb{R}^{d}\right)
$$

and $\mathcal{D}(A)=\left\{f \in L^{2}\left(\mathbb{R}^{d}\right): \int_{[0, \infty)} \gamma^{2} d\left(E_{\gamma} f, f\right)<\infty\right\}$ (see Fukushima et al. (1994) paragraph 1.3.4 p.18).

Consequently, the family of operators $\left\{T_{t} \stackrel{\text { def }}{=} \mathrm{e}^{t A}: t>0\right\}$ is a strongly continuous semigroup of self-adjoint contractions acting on $L^{2}\left(\mathbb{R}^{d}\right)$ (Fukushima et al. (1994) Lemma 1.3 .2 p.19) and

$$
\left(T_{t} f, g\right)=\int_{[0, \infty)} \mathrm{e}^{-\gamma t} d\left(E_{\gamma} f, g\right) \quad \forall f \in L^{2}\left(\mathbb{R}^{d}\right), g \in L^{2}\left(\mathbb{R}^{d}\right)
$$

Note that for any $\gamma \geq 0, t>0$, and any functions $f \in L^{2}\left(\mathbb{R}^{d}\right)$ and $g \in L^{2}\left(\mathbb{R}^{d}\right)$, we have the commutation property

$$
\begin{aligned}
\left(T_{t} E_{\gamma} f, g\right) & =\left(E_{\gamma} f, T_{t} g\right) \\
& =\int_{[0, \infty)} \mathrm{e}^{-\xi t} d_{\xi}\left(E_{\gamma} f, E_{\xi} g\right)=\int_{[0, \gamma]} \mathrm{e}^{-\xi t} d_{\xi}\left(E_{\xi} E_{\gamma} f, g\right)+\int_{[\gamma, \infty)} \mathrm{e}^{-\xi t} d_{\xi}\left(E_{\xi} E_{\gamma} f, g\right) \\
& =\int_{[0, \gamma]} \mathrm{e}^{-\xi t} d_{\xi}\left(E_{\xi} f, g\right)=\int_{[0, \infty)} \mathrm{e}^{-\xi t} d_{\xi}\left(E_{\xi} f, E_{\gamma} g\right)=\left(T_{t} f, E_{\gamma} g\right)=\left(E_{\gamma} T_{t} f, g\right) .
\end{aligned}
$$

Note also that for any $f \in L^{2}\left(\mathbb{R}^{d}\right)$ and any $t>0$,

$$
\begin{aligned}
\int_{[0, \infty)} \gamma^{2} d\left(E_{\gamma} T_{t} f, T_{t} f\right) & =\int_{[0, \infty)} \gamma^{2} d_{\gamma}\left(\int_{[0, \infty)} \mathrm{e}^{-\xi t} d_{\xi}\left(E_{\gamma} E_{\xi} f, T_{t} f\right)\right) \\
& =\int_{[0, \infty)} \gamma^{2} d_{\gamma}\left(\int_{[0, \infty)} \mathrm{e}^{-\xi t} d_{\xi}\left(\int_{[0, \infty)} \mathrm{e}^{-\theta t} d_{\theta}\left(E_{\gamma} E_{\xi} f, E_{\theta} f\right)\right)\right) \\
& =\int_{[0, \infty)} \gamma^{2} \mathrm{e}^{-2 \gamma t} d_{\gamma}\left(E_{\gamma} f, f\right) \\
& \leq \frac{4}{t^{2}} \mathrm{e}^{-2} \int_{[0, \infty)} \mathrm{e}^{-\gamma t} d_{\gamma}\left(E_{\gamma} f, f\right)=\frac{4}{t^{2}} \mathrm{e}^{-2}\left(T_{t} f, f\right) \leq \frac{4}{t^{2}} \mathrm{e}^{-2}\|f\|^{2}<+\infty,
\end{aligned}
$$

where we have used the spectral family property, the associativity of the Stieltjes integral and the inequality $\gamma^{2} \mathrm{e}^{-\gamma t} \leq 4 \mathrm{e}^{-2} / t^{2}$. The above inequality ensures that $T_{t} f \in \mathcal{D}(A)$ for any $t>0$.

From the fact that $\left|\frac{d}{d t} \mathrm{e}^{-\gamma t}\right| \leq \gamma$ is integrable w.r.t. $d\left(E_{\gamma} h, g\right)$ whenever $h \in \mathcal{D}(A)$, we deduce from the commutation property that for any $f, g \in L^{2}\left(\mathbb{R}^{d}\right)$ and for any $s>0$

$$
-\frac{d}{d t}\left(T_{t} f, T_{s} g\right)=\int_{[0, \infty)} \gamma \mathrm{e}^{-\gamma t} d\left(E_{\gamma} f, T_{s} g\right) \xrightarrow[t \searrow 0+]{\longrightarrow} \int_{[0, \infty)} \gamma d\left(E_{\gamma} T_{s} f, g\right)=\left(-A T_{s} f, g\right)
$$

where the limit exists and is well defined (since we have shown that $T_{s} f \in \mathcal{D}(A)$ ).
If moreover $g \in \mathcal{D}[\mathcal{E}]$ then

$$
\begin{align*}
-\frac{d}{d s}\left(T_{s} f, g\right) & =-\left.\frac{d}{d t}\left(T_{s+t} f, g\right)\right|_{t=0+}=-\left.\frac{d}{d t}\left(T_{t} f, T_{s} g\right)\right|_{t=0+} \\
& =\left(-A T_{s} f, g\right)=\mathcal{E}\left(T_{s} f, g\right) \tag{2.2}
\end{align*}
$$

And since $-\frac{d}{d s}\left(T_{s} f, g\right)=-\frac{d}{d s}\left(T_{s} g, f\right)$ by the symmetry property of $T_{s}$, we deduce that $\mathcal{E}\left(T_{s} f, g\right)=$ $\mathcal{E}\left(f, T_{s} g\right)$ for any $f, g \in \mathcal{D}[\mathcal{E}]$.

Consequently, for any $f \in \mathcal{D}[\mathcal{E}]$ and using the ellipticity of the coefficient matrix $a$,

$$
\begin{aligned}
\lambda\left\|\nabla T_{s} f\right\|^{2} \leq & \mathcal{E}\left(T_{s} f, T_{s} f\right)=\mathcal{E}\left(T_{2 s} f, f\right)=\left(-A T_{2 s} f, f\right) \\
& =\int_{[0, \infty)} \gamma \mathrm{e}^{-2 \gamma s} d\left(E_{\gamma} f, f\right) \leq \frac{\mathrm{e}^{-1}}{s} \int_{[0, \infty)} \mathrm{e}^{-\gamma s} d\left(E_{\gamma} f, f\right)=\frac{\mathrm{e}^{-1}}{s}\left(T_{s} f, f\right) \leq \frac{\|f\|^{2}}{s}
\end{aligned}
$$

from which we deduce the fundamental estimate

$$
\begin{equation*}
\left\|\nabla T_{s} f\right\| \leq \frac{\|f\|}{\sqrt{\lambda s}}, \quad \forall s>0 \tag{2.3}
\end{equation*}
$$

In turn this estimate implies that for any $f \in L^{2}\left(\mathbb{R}^{d}\right), g \in \mathcal{D}[\mathcal{E}]$, the function

$$
s \mapsto \mathcal{E}\left(T_{s} f, g\right) \text { is integrable on }(0, t],
$$

and from (2.2) and the right continuity of $s \mapsto T_{s} f$ at time $s=0+$ (one may extend $T_{0} f=f$ as long as no differentiation of $s \mapsto T_{s} f$ is implied at $s=0+$ when $\left.f \notin \mathcal{D}(A)\right)$, we deduce the integrated version of (2.2) namely
$\forall f \in L^{2}\left(\mathbb{R}^{d}\right), \forall g \in \mathcal{D}[\mathcal{E}]$,

$$
\begin{equation*}
\left(T_{t} f, g\right)-(f, g)=-\int_{0}^{t} \sum_{i, j=1}^{d}\left(a_{i j} D_{j} T_{s} f, D_{i} g\right) d s=-\int_{0}^{t} \mathcal{E}\left(T_{s} f, g\right) d s, \quad t \in(0, \infty) \tag{2.4}
\end{equation*}
$$

2.2. Link with the results of D.W. Stroock Stroock (1988). In his celebrated article Diffusion semigroups corresponding to uniformly elliptic divergence form operators D.W. Stroock constructs via a regularization procedure a Feller continuous semigroup $\left\{P_{t}: t>0\right\}$ associated to $a$ with the properties that (with our notations)
(1) the map $t \in[0, \infty) \mapsto P_{t} \phi \in H^{1}\left(\mathbb{R}^{d}\right)$ is a weakly continuous map for each $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$.
(2) $\forall \phi, \psi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$,

$$
\begin{equation*}
\left(P_{t} \phi, \psi\right)-(\phi, \psi)=-\int_{0}^{t}\left(a \nabla P_{s} \phi, \nabla \psi\right) d s=-\int_{0}^{t} \mathcal{E}\left(P_{s} \phi, \psi\right) d s, \quad t \in(0, \infty) . \tag{2.5}
\end{equation*}
$$

(Nota: please note that there is a sign error in the original version of Stroock (1988)).
In fact, $\left\{P_{t}: t>0\right\}$ determines a unique strongly continuous semigroup $\left\{\bar{P}_{t}: t>0\right\}$ of self-adjoint contractions on $L^{2}\left(\mathbb{R}^{d}\right)$.

The aim of this subsection is to prove the following equality, which clarifies the relationship between the results obtained in Stroock (1988) and the those provided by the theory of Dirichlet forms Fukushima et al. (1994).

## Proposition 2.1.

$$
\begin{equation*}
\left\{\bar{P}_{t}: t>0\right\}=\left\{T_{t}: t>0\right\} \quad \text { on } L^{2}\left(\mathbb{R}^{d}\right) . \tag{2.6}
\end{equation*}
$$

Proof: The semigroup $\left\{\bar{P}_{t}: t>0\right\}$ is strongly continuous on $H^{1}\left(\mathbb{R}^{d}\right)$. Moreover, for each $t>0$, the operator $\bar{P}_{t}$ maps $L^{2}\left(\mathbb{R}^{d}\right)$ into $H^{1}\left(\mathbb{R}^{d}\right)$ and for each $f \in H^{1}\left(\mathbb{R}^{d}\right)=\mathcal{D}[\mathcal{E}]$, we have the fundamental estimate

$$
\begin{equation*}
\left\|\nabla \bar{P}_{s} f\right\| \leq \frac{1}{\sqrt{\lambda}}\left(\frac{\|f\|}{\sqrt{s}}\right) \wedge\|\nabla f\|, \quad \forall s>0 . \tag{2.7}
\end{equation*}
$$

(See Stroock (1988) Theorem II.3.1. p.341).
This estimate implies that for for each $f, g \in \mathcal{D}[\mathcal{E}]$ and any $t, s>0$,

$$
\left|\mathcal{E}\left(\bar{P}_{t} f, g\right)-\mathcal{E}\left(\bar{P}_{s} f, g\right)\right| \leq \Lambda\|\nabla g\| \frac{\left\|\bar{P}_{t \vee s-t \wedge s} f-f\right\|}{\sqrt{\lambda(t \wedge s)}} \underset{s \rightarrow t}{\longrightarrow} 0
$$

which ensures the continuity of $s \mapsto \mathcal{E}\left(\bar{P}_{s} \phi, \psi\right)$ for any $\phi, \psi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. Since $\left(\bar{P}_{t}\right)_{t>0}$ and $\left(P_{t}\right)_{t>0}$ coïncide on $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, we may differentiate in (2.5) (as long as $t>0$ ) to find that

$$
\begin{equation*}
\frac{d}{d t}\left(\bar{P}_{t} \phi, \psi\right)=-\mathcal{E}\left(\bar{P}_{t} \phi, \psi\right), \quad t \in(0, \infty) . \tag{2.8}
\end{equation*}
$$

This has to be compared to (2.2).
Let us now justify rigorously that for any $t>0, s \in(0, t)$ and $\phi, \psi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$,

$$
\begin{equation*}
\frac{d}{d s}\left(T_{s} \phi, \bar{P}_{t-s} \psi\right)=\left.\frac{d}{d u}\left(T_{u} \phi, \bar{P}_{t-s} \psi\right)\right|_{u=s}-\left.\frac{d}{d u}\left(T_{s} \phi, \bar{P}_{t-u} \psi\right)\right|_{u=s} \tag{2.9}
\end{equation*}
$$

We have for sufficiently small $0 \leq h<t-s$ and using the strong continuity of $\left(\bar{P}_{t}\right)_{t>0}$

$$
\begin{aligned}
\left|\left(T_{s+h} \phi-T_{s} \phi, \bar{P}_{t-s+h} \psi-\bar{P}_{t-s} \psi\right)\right| & \leq\left\|\bar{P}_{t-s+h} \psi-\bar{P}_{t-s} \psi\right\|\left\|T_{s+h} \phi-T_{s} \phi\right\| \\
& \leq \varepsilon_{\psi}(h)\left(\int_{[0, \infty)} \mathrm{e}^{-2 \gamma s}\left(\mathrm{e}^{-\gamma h}-1\right)^{2} d\left(E_{\gamma} \phi, \phi\right)\right)^{1 / 2} \\
& \leq \varepsilon_{\psi}(h)\left(\int_{[0, \infty)} \mathrm{e}^{-2 \gamma s}(\gamma h)^{2} d\left(E_{\gamma} \phi, \phi\right)\right)^{1 / 2} \\
& \leq h \varepsilon_{\psi}(h)\left(\int_{[0, \infty)} \mathrm{e}^{-\gamma s}\left(\gamma^{2} \mathrm{e}^{-\gamma s}\right) d\left(E_{\gamma} \phi, \phi\right)\right)^{1 / 2} \\
& \leq h \varepsilon_{\psi}(h) \frac{2 \mathrm{e}^{-1}}{s}\|\phi\|
\end{aligned}
$$

where as usual $\varepsilon_{\psi}(\cdot)$ denotes some positive continuous function vanishing at zero. We deduce that

$$
\frac{1}{h}\left|\left(T_{s+h} \phi-T_{s} \phi, \bar{P}_{t-s+h} \psi-\bar{P}_{t-s} \psi\right)\right| \underset{h \rightarrow 0}{\longrightarrow} 0
$$

implying (2.9).
Hence, from (2.9) and applying (2.2) and (2.8), we have that

$$
\begin{equation*}
\frac{d}{d s}\left(T_{s} \phi, \bar{P}_{t-s} \psi\right)=0, \quad s \in(0, t) \tag{2.10}
\end{equation*}
$$

Integrating the identity (2.10) on $(0, t)$ and using the time continuity of both semigroups $\left(T_{t}\right)$ and $\left(\bar{P}_{t}\right)$ up to time $s=0+$ gives

$$
\begin{equation*}
\left(T_{t} \phi, \psi\right)=\left(\phi, \bar{P}_{t} \psi\right)=\left(\bar{P}_{t} \phi, \psi\right) \tag{2.11}
\end{equation*}
$$

which holds for any $\phi, \psi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. Since $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ is dense in $L^{2}\left(\mathbb{R}^{d}\right)$, using the strong continuity of both semigroups $\left(T_{t}\right)$ and $\left(\bar{P}_{t}\right)$, we finally deduce from (2.11) the identification (2.6).

Consequently, all results in Stroock (1988) that are valid for $\left\{\bar{P}_{t}: t>0\right\}$ are true for $\left\{T_{t}: t>0\right\}$. For example, identifying abusively $\left\{T_{t}: t>0\right\}$ with its Feller restriction $\left\{P_{t}: t>0\right\}$ on $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, we deduce that there is a $p \in C\left((0, \infty) \times \mathbb{R}^{d} \times \mathbb{R}^{d}\right)$ such that

$$
\begin{equation*}
\left[T_{t} \phi\right](x)=\int_{\mathbb{R}^{d}} \phi(y) p(t, x, y) d y, \quad \ell(d x)-\text { a.e., } \quad \phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right) \tag{2.12}
\end{equation*}
$$

Moreover, the fundamental function $p$ satisfies the well-known Aronson's estimates for the fundamental solutions of elliptic divergence form operators, namely there exists a constant $M(\lambda, \Lambda, d) \in$ $[1, \infty)$ such that

$$
\begin{equation*}
\frac{1}{M t^{d / 2}} \exp \left(-M|x-y|^{2} / t\right) \leq p(t, x, y) \leq \frac{M}{t^{d / 2}} \exp \left(-|x-y|^{2} / M t\right) \tag{2.13}
\end{equation*}
$$

Finally, we have the convergence result of Stroock (1988) (Theorem II.3.1. p.341) that we state roughly without introducing the necessary notations (see Stroock (1988) for details): if $\left\{a_{n}\right\}_{1}^{\infty} \subset$ $\mathcal{A}(\lambda, \Lambda)$ and $a_{n} \longrightarrow a$ almost everywhere, then $p^{n}(t, x, y) \longrightarrow p(t, x, y)$ uniformly on compacts (in $\left.(0, \infty) \times \mathbb{R}^{d} \times \mathbb{R}^{d}\right)$ and for each $t \in[0, \infty)$ and $\phi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right), T_{t}^{n} \phi \longrightarrow T_{t} \phi$ in $H^{1}\left(\mathbb{R}^{d}\right)$.

## 3. Associated stochastic processes

Since $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ is a regular Dirichlet form (with the space $\mathcal{D}[\mathcal{E}] \cap C_{c}\left(\mathbb{R}^{d}\right)$ or $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ as a special standard core, see e.g. Exercice 1.4.1 in Fukushima et al. (1994)), we are in position to apply Theorem 7.2.1 p. 380 of Fukushima et al. (1994).

We may associate to $(\mathcal{E}, \mathcal{D}[\mathcal{E})]$ and its corresponding semigroup $\left(T_{t}\right)$ a Hunt process, symmetric w.r.t the Lebesgue measure $\ell(d x)$ on $\mathbb{R}^{d}$. We shall denote by $\mathbb{M}=\left(\Omega,\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathcal{F},\left(X_{t}\right)_{t \geq 0},\left(P^{x}\right)_{x \in \mathbb{R}^{d}}\right)$ this Hunt process, with $X=\left(X^{1}, \ldots, X^{d}\right)$. The correspondence with $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ and $\left(\bar{T}_{t}\right)$ is through

$$
\begin{equation*}
E^{x}\left[f\left(X_{t}\right)\right]=T_{t} f(x), \quad \forall f \in L^{2}\left(\mathbb{R}^{d}\right), \forall t \geq 0, \forall x \in \mathbb{R}^{d} \tag{3.1}
\end{equation*}
$$

(see the discussion p160 in Fukushima et al. (1994), at the beginning of Section 4.2).
The aim of this section is to present various representations of $\mathbb{M}$ in various contexts. We start with the general case and then present a Skorokhod decomposition in the specific case where our Dirichlet form is associated to a transmission parabolic operator in divergence form.
3.1. The Revuz correspondence for additive functionals and the Fukushima decomposition.
3.1.1. A reminder of the Revuz correspondence for additive functionals. Denote by $\left\{R_{\alpha}: \alpha>0\right\}$ the Markovian resolvent kernel of the Markovian transition function $\{\hat{p}(t, x, d y):=p(t, x, y) d y: t>0\}$. Then, for any $\alpha>0, f \in \mathcal{B}_{b}\left(\mathbb{R}^{d}\right)$ and $x \in \mathbb{R}^{d}, R_{\alpha} f(x)=\int_{\mathbb{R}^{d}} r_{\alpha}(x, y) f(y) d y$ with $r_{\alpha}(x, y)=$ $\int_{0}^{\infty} \mathrm{e}^{-\alpha t} p(t, x, y) d t$.

Denote by $S$ the set of positive Radon measures on $\left(\mathbb{R}^{d}, \mathcal{B}\left(\mathbb{R}^{d}\right)\right.$ ). For $\mu \in S$ define $R_{1} \mu(x)=$ $\int_{\mathbb{R}^{d}} r_{1}(x, y) \mu(d y)\left(x \in \mathbb{R}^{d}\right)$ and introduce the subset of finite energy measures

$$
S_{0}:=\left\{\mu \in S: \sup _{v \in \mathcal{D}[\mathcal{E}] \cap C_{c}\left(\mathbb{R}^{d}\right)} \int_{\mathbb{R}^{d}} \frac{|v(x)|}{\|v\|_{\mathcal{E}_{1}}} \mu(d x)<\infty\right\},
$$

(where we follow the notations of Fukushima et al. (1994)).
Finally, introduce

$$
S_{00}:=\left\{\mu \in S_{0}: \mu\left(\mathbb{R}^{d}\right)<\infty,\left\|R_{1} \mu(.)\right\|_{\infty}<\infty\right\}
$$

Let us denote respectively by $\mathbf{A}_{c}^{+}$and $\mathbf{A}_{c, 1}^{+}$the families of all Positive Continuous Additive Functionals (PCAF in short) (resp. the family of all PCAF in the strict sense) associated to $\mathbb{M}$ (for the distinction between $\mathbf{A}_{c}^{+}$and $\mathbf{A}_{c, 1}^{+}$, see Fukushima et al. (1994) the introduction of Section 5.1).

The Revuz correspondence asserts that there is a one-to-one correspondence (up to equivalence of processes) between $\mathbf{A}_{c}^{+}$and $S$. This correspondence permits to construct for any $\mu \in S_{00}$ a unique PCAF in the strict sense $A \in \mathbf{A}_{c, 1}^{+}$such that

$$
\begin{equation*}
\forall x \in \mathbb{R}^{d}, \quad E^{x} \int_{0}^{\infty} \mathrm{e}^{-t} d A_{t}=R_{1} \mu(x) \tag{3.2}
\end{equation*}
$$

(see for e.g. Theorem 5.1.4 in Fukushima et al. (1994)).
In order to get a bijective map, introduce a new subset $S_{1}$ of $S$ defined by $\mu \in S_{1}$ if there exists a sequence $\left(E_{n}\right)_{n \geq 0}$ of Borel finely open sets increasing to $\mathbb{R}^{d}$ satisfying that $\mathbb{I}_{E_{n}} \cdot \mu \in S_{00}$ for each $n$. Then, there is a one-to-one correspondence between $S_{1}$ and $\mathbf{A}_{c, 1}^{+}$(up to equivalence) which is given by relation (3.2) whenever $\mu \in S_{00}$. The set of measures $S_{1}$ is called the set of smooth measures (in the strict sense).

Let us introduce $\mathcal{D}[\mathcal{E}]_{b}$ (resp. $\mathcal{D}[\mathcal{E}]_{b, \text { loc }}$ ) the space of essentially bounded functions belonging to $\mathcal{D}[\mathcal{E}]$ (resp. locally to $\mathcal{D}[\mathcal{E}]$. A function $u$ is in $\mathcal{D}[\mathcal{E}]_{b, \text { loc }}$ if for any compact set $G$, there exists a bounded function $\omega \in \mathcal{D}[\mathcal{E}]$ such that $u=\omega, \quad \ell(d x)$-a.e. on $G)$.

For $u \in \mathcal{D}[\mathcal{E}]_{b}$, we may associate a unique positive Radon measure $\mu_{\langle u\rangle} \in S$, satisfying

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} f(x) \mu_{\langle u\rangle}(d x)=2 \mathcal{E}(u f, u)-\mathcal{E}\left(u^{2}, f\right), \quad \forall f \in \mathcal{D}[\mathcal{E}] \cap C_{c}\left(\mathbb{R}^{d}\right) \tag{3.3}
\end{equation*}
$$

Observe that the positivity of the right hand side in (3.3) comes from

$$
\begin{aligned}
0 \leq E^{x}\left[\left(u\left(X_{t}\right)-u(x)\right)^{2}\right] & =E^{x}\left[u^{2}\left(X_{t}\right)-2 u(x) u\left(X_{t}\right)+u^{2}(x)\right] \\
& =\left(T_{t} u^{2}-u^{2}\right)(x)-2 u(x)\left(T_{t} u-u\right)(x)
\end{aligned}
$$

Taking the scalar product with $f \in \mathcal{D}[\mathcal{E}] \cap C_{c}\left(\mathbb{R}^{d}\right)$ and dividing by $t$ gives

$$
0 \leq\left(\frac{T_{t} u^{2}-u^{2}}{t}, f\right)-2\left(\frac{T_{t} u-u}{t}, u f\right)
$$

which tends to $2 \mathcal{E}(u f, u)-\mathcal{E}\left(u^{2}, f\right)$ as $t$ tends to $0+$ whenever $u \in \mathcal{D}(A)$; the positivity in the case $u \in \mathcal{D}[\mathcal{E}]_{b}$ may be obtained by a density argument.

If $u \in \mathcal{D}[\mathcal{E}]_{b, \text { loc }}$, we may construct $\mu_{\langle u\rangle} \in S$ with the help of a sequence $\left(G_{n}\right)_{n \geq 0}$ of relatively compact open sets such that $\overline{G_{n}} \subset G_{n+1}$ and $\bigcup_{n \geq 0} G_{n}=\mathbb{R}^{d}$. Let $\left(u_{n}\right)_{n \geq 0}$ a sequence of functions in $\mathcal{D}[\mathcal{E}]_{b}$ satisfying $u_{n}=u$ on $G_{n}$. There is no ambiguïty in defining $\mu_{\langle u\rangle}=\mu_{\left\langle u_{n}\right\rangle}$ on $G_{n}$ because the construction is consistent (since $\mu_{\left\langle u_{n}\right\rangle}=\mu_{\left\langle u_{n+1}\right\rangle}$ on $G_{n}$ ). For an account on the above assertions, please refer to Fukushima et al. (1994) Section 3.2.
3.1.2. The Fukushima decomposition. Note that obviously $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ is strong local, so we may apply Theorem 5.5.5 in Fukushima et al. (1994).

Suppose that a function $u$ satisfies the following conditions:
(1) $u \in \mathcal{D}[\mathcal{E}]_{b, \text { loc }}, u$ is finely continuous on $\mathbb{R}^{d}$.
(2) $\mathbb{I}_{G} \cdot \mu_{\langle u\rangle} \in S_{00}$ for any relatively compact open set $G$.
(3) $\exists \varrho=\varrho^{(1)}-\varrho^{(2)}$ with $\mathbb{I}_{G} \cdot \varrho^{(1)}, \mathbb{I}_{G} \cdot \varrho^{(2)} \in S_{00}$ for any relatively compact open set $G$ and

$$
\mathcal{E}(u, v)=(\varrho, v), \quad \forall v \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)
$$

(Note that even though $u$ is not formally in $\mathcal{D}[\mathcal{E}]$, the quantity $\mathcal{E}(u, v)$ is well-defined because $v$ has compact support and $\left.u \in \mathcal{D}[\mathcal{E}]_{b, \text { loc }}\right)$.
Let $A^{(1)}, A^{(2)}$, and $B$ be PCAF's in the strict sense with Revuz measures $\varrho^{(1)}, \varrho^{(2)}$, and $\mu_{\langle u\rangle}$ respectively. Then, Theorem 5.5.5 in Fukushima et al. (1994) asserts that

$$
\begin{equation*}
u\left(X_{t}\right)-u\left(X_{0}\right)=M_{t}^{[u]}+N_{t}^{[u]}, \quad \mathbb{P}^{x}-a . s, \forall x \in \mathbb{R}^{d} \tag{3.4}
\end{equation*}
$$

Here,

$$
\begin{equation*}
N^{[u]}=-A^{(1)}+A^{(2)}, \quad \mathbb{P}^{x}-a . s, \forall x \in \mathbb{R}^{d} \tag{3.5}
\end{equation*}
$$

and $M^{[u]}$ is a local Additive Functional in the strict sense such that for any relatively compact set $G$,

$$
E^{x} M_{t \wedge \tau_{G}}^{[u]}=0, \quad \forall x \in G
$$

and

$$
E^{x}\left[\left(M_{t \wedge \tau_{G}}^{[u]}\right)^{2}\right]=\mathbb{E}^{x} B_{t \wedge \tau_{G}}, \quad \forall x \in G
$$

where $\tau_{G}=\inf \left(s>0: X_{s} \notin G\right)$ stands for the first leaving time from $G$ (with the convention $\inf \emptyset=\infty)$ and $B$ denotes the PCAF in the strict sense with Revuz measure $\mu_{\langle u\rangle}$.
3.2. An insight to the Lyons-Zheng decomposition for diffusions associated to divergence form operators. For the seek of completeness, in this paragraph we present briefly an insight on the LyonsZheng decomposition Lyons and Zheng (1990) for the stochastic process in correspondence with a divergence form operator.
3.2.1. Time reversal for diffusions associated to divergence form operators. Assume for a moment that $a \in \mathcal{A}(\lambda, \Lambda)$ is very smooth and belongs to $C^{\infty}\left(\mathbb{R}^{d} \rightarrow \mathcal{M}_{d}(\mathbb{R})\right)$.

In this smooth case, the process $\left\{X_{t}=\left(X_{t}^{1}, \ldots, X_{t}^{d}\right): 0 \leq t \leq T\right\}$ in correspondence with $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ and constructed in the previous section becomes a diffusion process with values in $\mathbb{R}^{d}$ and well-known results from the general theory for solutions of stochastic differential equations ensure that $\left(X_{t}\right)$ is solution of

$$
\begin{equation*}
X_{t}^{k}=x_{k}+\int_{0}^{t} \sum_{j=1}^{d} \sigma_{k j}\left(X_{s}\right) d W_{s}^{j}+\int_{0}^{t} \sum_{j=1}^{d} \partial_{j} a_{k j}\left(X_{s}\right) d s \quad t \geq 0, \mathbb{P}^{x}-a . s ., \forall x=\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d} \tag{3.6}
\end{equation*}
$$

where $W=\left(W^{1}, \ldots, W^{d}\right)$ denotes a $d$-dimensional standard Brownian motion starting from zero and $\sigma: \mathbb{R}^{d} \rightarrow \mathcal{M}_{d}(\mathbb{R})$ denotes the positive square-root of coefficient $2 a$ i.e. the positive matrix real valued coefficient satisfying

$$
\sigma \sigma^{*}(x)=2 a(x), \quad \forall x \in \mathbb{R}^{d}
$$

We will denote by $\mathcal{L}$ the classical generator of $\left(X_{t}\right)$,

$$
\begin{equation*}
\mathcal{L}=\sum_{i, j=1}^{d} \partial_{i}\left(a_{i j} \partial_{j}\right)=\sum_{i, j=1}^{d} a_{i j} \partial_{i j}^{2}+\sum_{i, j=1}^{d}\left(\partial_{i} a_{i j}\right) \partial_{j} \tag{3.7}
\end{equation*}
$$

acting on $C^{2}\left(\mathbb{R}^{d}\right)$ real valued functions.
Fix $x=\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d}$. For any fixed $0 \leq t<1$, denote $\overline{\mathcal{L}}_{t}^{x}$ the second order differential operator

$$
\overline{\mathcal{L}}_{t}^{x}=\sum_{i, j=1} a_{i j} \partial_{i j}^{2}-\sum_{i, j=1}^{d}\left(\partial_{i} a_{i j}\right) \partial_{j}+[p(1-t, x, .)]^{-1} \sum_{i, j=1}^{d} \partial_{j}\left(\left(2 a_{i j}\right) p(1-t, x, .)\right) \partial_{i}
$$

acting also on $C^{2}\left(\mathbb{R}^{d}\right)$ real valued functions and where $p$ denotes the fundamental solution (2.12). (Nota: the reader should be careful with the 2 factor appearing in the last integral, that is due to the fact that we are considering $\sigma$ as the square root of $2 a$ and not $a$. Also, since $a$ is assumed to be smooth, classical estimates for $p$ ensure that the term $p(1-t, x, .)^{-1} \partial_{j}\left(\left(2 a_{i j}\right) p(1-t, x,).\right) \partial_{i}$ is well defined).

Consider $\left\{\bar{X}_{t}:=X_{1-t}: t \in[0,1)\right\}$ the time reversed process of $\left(X_{t}\right)$. The fact that the time reversal of a Markov process is again a (weak) Markov process with respect to the reversed filtration traces back to the seminal result of K.L. Chung and J.B. Walsh Chung and Walsh (1969). Moreover, it is proved in Haussmann and Pardoux (1986) that $\left(\bar{X}_{t}\right)_{0 \leq t<1}$ is a diffusion process with generator $\overline{\mathcal{L}}$ up to time 1 excluded. The process $\left\{\bar{X}_{t}=\left(X_{1-t}^{1}, \ldots, X_{1-t}^{d}\right): 0 \leq t<1\right\}$ is solution of

$$
\begin{aligned}
\bar{X}_{t}^{k}=X_{1} & +\int_{0}^{t} \sum_{j=1}^{d} \sigma_{k j}\left(\bar{X}_{s}\right) d \beta_{s}^{j}-\int_{0}^{t} \sum_{j=1}^{d} \partial_{j} a_{k j}\left(\bar{X}_{s}\right) d s \\
& +\int_{0}^{t}\left[p\left(1-t, x, \bar{X}_{s}\right)\right]^{-1} \sum_{i, j=1}^{d} \partial_{j}\left(\left(2 a_{i j}\right) p\left(1-s, x, \bar{X}_{s}\right)\right) d s, \mathbb{P}^{x}-\text { a.s., } \forall x=\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d}
\end{aligned}
$$

where $\beta=\left(\beta^{1}, \ldots, \beta^{d}\right)$ denotes a $d$-dimensional standard Brownian motion starting from zero and adapted to the filtration $\left\{\mathcal{F}_{t}^{\bar{X}}:=\sigma\left(X_{1-u}: u \in[0, t]\right) \mid t \in[0,1)\right\}$.
3.2.2. The Lyons-Zheng decomposition. We now make the following observation. Fix $0<\varepsilon<1$, then for any $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and arbitrary $t \in[\varepsilon, 1]$

$$
\begin{aligned}
\varphi\left(X_{t}\right)-\varphi(x) & =\frac{1}{2}\left(\varphi\left(X_{t}\right)-\varphi(x)\right)+\frac{1}{2}\left[\left(\varphi\left(X_{1-(1-t)}\right)-\varphi\left(X_{1-(1-\varepsilon)}\right)\right)+\left(\varphi\left(X_{\varepsilon}\right)-\varphi(x)\right)\right] \\
& =\frac{1}{2}\left(\varphi\left(X_{t}\right)-\varphi(x)\right)-\frac{1}{2}\left(\varphi\left(\bar{X}_{1-\varepsilon}\right)-\varphi\left(\bar{X}_{1-t}\right)\right)+\frac{1}{2}\left(\varphi\left(X_{\varepsilon}\right)-\varphi(x)\right)
\end{aligned}
$$

and applying Itô's formula, we find that for all $t \in[\varepsilon, 1]$

$$
\begin{align*}
& \varphi\left(X_{t}\right)-\varphi(x) \\
&= \frac{1}{2}\left(\int_{0}^{t}\left(\nabla \varphi\left(X_{u}\right), \sigma\left(X_{u}\right) d W_{u}\right)+\int_{0}^{t} \mathcal{L} \varphi\left(X_{u}\right) d u\right)-\frac{1}{2}\left(\int_{1-t}^{1-\varepsilon}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right)-\int_{1-t}^{1-\varepsilon} \overline{\mathcal{L}}_{u}^{x} \varphi\left(\bar{X}_{u}\right) d u\right) \\
&+\frac{1}{2}\left(\varphi\left(X_{\varepsilon}\right)-\varphi(x)\right) \\
&= \frac{1}{2} M_{t}^{\varphi}+\frac{1}{2} \tilde{N}_{t}^{\varphi, \varepsilon}+\frac{1}{2} \int_{\varepsilon}^{t}\left(\mathcal{L}-\overline{\mathcal{L}}_{1-u}^{x}\right) \varphi\left(X_{u}\right) d u+\frac{1}{2}\left(\varphi\left(X_{\varepsilon}\right)-\varphi(x)+\int_{0}^{\varepsilon} \mathcal{L} \varphi\left(X_{u}\right) d u\right) \tag{3.8}
\end{align*}
$$

where we have set $M_{t}^{\varphi}:=\int_{0}^{t}\left(\nabla \varphi\left(X_{u}\right), \sigma\left(X_{u}\right) d W_{u}\right)$ and $\tilde{N}_{t}^{\varphi, \varepsilon}:=-\int_{1-t}^{1-\varepsilon}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right)$.
We may write

$$
\begin{aligned}
\tilde{N}_{t}^{\varphi, \varepsilon} & =-\int_{1-t}^{1-\varepsilon}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right) \\
& =\int_{0}^{1-t}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right)-\int_{0}^{1-\varepsilon}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right) \\
& =N_{1-t}^{\varphi}-N_{1-\varepsilon}^{\varphi},
\end{aligned}
$$

where $\left(N_{\theta}^{\varphi}\right)$ stands for the $\mathcal{F}^{\bar{X}}$ martingale $\left(\int_{0}^{\theta}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right)\right)_{\theta \in[0,1)}$.
Now for any $z \in \mathbb{R}^{d}$

$$
\begin{aligned}
\frac{1}{2}\left(\mathcal{L}-\overline{\mathcal{L}}_{1-u}^{x}\right) \varphi(z) & =\sum_{i, j=1}^{d}\left(\partial_{i} a_{i j}\right) \partial_{j} \varphi(z)-[p(u, x, z)]^{-1} \sum_{i, j=1}^{d} \partial_{j}\left(a_{i j} p(u, x, z)\right) \partial_{i} \varphi(z) \\
& =\left[\sum_{i, j=1}^{d}\left(\partial_{i} a_{i j}\right) \partial_{j} \varphi(z)-\sum_{i, j=1}^{d}\left(\partial_{j} a_{i j}\right) \partial_{i} \varphi(z)\right]-[p(u, x, z)]^{-1} \sum_{i, j=1}^{d} a_{i j} \partial_{j} p(u, x, z) \partial_{i} \varphi(z) \\
& =-[p(u, x, z)]^{-1} \sum_{i, j=1}^{d} a_{i j} \partial_{j} p(u, x, z) \partial_{i} \varphi(z)
\end{aligned}
$$

where terms cancel due to the symmetry of the coefficient matrix $a$.
Of course $\lim _{\varepsilon \searrow 0+}\left(\varphi\left(X_{\varepsilon}\right)-\varphi(x)+\int_{0}^{\varepsilon} \mathcal{L} \varphi\left(X_{s}\right) d s\right)=0$. Applying the martingale convergence theorem ensures that the martingale term $N_{1-\varepsilon}^{\varphi}$ tends $\mathbb{P}^{x}-$ a.s. to $N_{1}^{\varphi}=\int_{0}^{1}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right)$ as $\varepsilon$ tends to zero.

In turn, this ensures from (3.8) that the limit

$$
\begin{equation*}
\lim _{\varepsilon \searrow 0+} \int_{\varepsilon}^{t}\left[p\left(u, x, X_{u}\right)\right]^{-1} \sum_{i, j=1}^{d} a_{i j} \partial_{j} p\left(u, x, X_{u}\right) \partial_{i} \varphi\left(X_{u}\right) d u \tag{3.9}
\end{equation*}
$$

exists $\mathbb{P}^{x}-$ a.s.
Coming back to (3.8) and taking limits as $\varepsilon \searrow 0+$ yields the Lyons-Zheng decomposition of $\left(X_{t}\right)$ when the coefficient matrix $a$ is smooth: for all $\varphi \in C_{c}^{\infty}, \forall t \in[0,1], \mathbb{P}^{x}-a . s$.

$$
\begin{equation*}
\varphi\left(X_{t}\right)-\varphi(x)=\frac{1}{2} M_{t}^{\varphi}+\frac{1}{2} \tilde{N}_{t}^{\varphi}-\int_{0}^{t}\left[p\left(u, x, X_{u}\right)\right]^{-1} \sum_{i, j=1}^{d} a_{i j} \partial_{j} p\left(u, x, X_{u}\right) \partial_{i} \varphi\left(X_{u}\right) d u \tag{3.10}
\end{equation*}
$$

with $\tilde{N}_{t}^{\varphi}=N_{1-t}^{\varphi}-N_{1}^{\varphi}$ an increment of a time-reversed martingale. The quadratic variations of the martingales involved in the decomposition are given by

$$
\left\langle M^{\varphi}\right\rangle_{t}=\int_{0}^{t} \sum_{i, j=1}^{d} 2 a_{i j}\left(X_{s}\right) \partial_{i} \varphi\left(X_{s}\right) \partial_{j} \varphi\left(X_{s}\right) d s, \text { and }\left\langle N^{\varphi}\right\rangle_{t}=\int_{0}^{t} \sum_{i, j=1}^{d} 2 a_{i j}\left(\bar{X}_{s}\right) \partial_{i} \varphi\left(\bar{X}_{s}\right) \partial_{j} \varphi\left(\bar{X}_{s}\right) d s
$$

We make the following observations:

- If we have proved (3.9), it is important to notice however that there is no limit to the deterministic quantity $\int_{\varepsilon}^{t}[p(u, x, z)]^{-1} \sum_{i, j=1}^{d} a_{i j} \partial_{j} p(u, x, z) \partial_{i} \varphi(z) d u$ as $\varepsilon$ tends to $0+$. This can be easily seen by performing the computation from the explicit Laplacian case where $p$ is just the Gaussian transition density of some Browian motion started from $x$.
- For any $\varepsilon>0$ the random variable $\tilde{N}_{t}^{\varphi, \varepsilon}=-\int_{1-t}^{1-\varepsilon}\left(\nabla \varphi\left(\bar{X}_{u}\right), \sigma\left(\bar{X}_{u}\right) d \beta_{u}\right)$ is measurable w.r.t the sigma field

$$
\sigma\left(\bar{X}_{u}: u \in[1-t, 1-\varepsilon]\right)=\sigma\left(X_{u}: u \in[\varepsilon, t]\right)
$$

so that $\left(\tilde{N}_{t}^{\varphi}\right)$ and all terms in (3.10) are adapted to $\left(\mathcal{F}_{t}^{X}\right)_{t \in[0,1]}$ the natural filtration of $X$.

- Only $\varphi$ and its first order partial derivatives appear in equality (3.10). Using a density argument and a little work we may prove that the equality holds for $\varphi \in H^{1}\left(\mathbb{R}^{d}\right)$.
- None of the quantities in the right hand side of (3.10) involve the derivatives of the coefficient matrix $a$ : the dependence on the derivatives of the coefficient $a$ is totally encompassed in the logarithmic derivative of the fundamental solution $p(t, x, y)$.
The idea is now to pick $a$ measurable in $\mathcal{A}(\lambda, \Lambda)$ and to take a sequence of smooth $\left\{a_{n}\right\}_{1}^{\infty} \subset$ $\mathcal{A}(\lambda, \Lambda)$ such that $a_{n} \longrightarrow a$ almost everywhere and to prove that there is convergence in law for the decomposition (3.10) for any $\varphi$ belonging to a the widdest possible class of functions. This programm has been successively performed in Rozkosz (1996a) and Rozkosz (1996b) in a more general setting of an inhomogeneous divergence operator.

Though theoretically powerful, the Lyons-Zheng decomposition is unlikely to be directly exploitable from a numerical perspective, as we do not have access to the logarithm derivative of the transition density, even in mild cases.

## 4. Stochastic dynamics associated to transmission operators in divergence form

4.1. Skorokhod representation of the Hunt process associated to a transmission operator in divergence form. Consider $\mathbb{R}^{d}=\bar{D}_{+} \cup D_{-}$with $D_{+}$and $D_{-}$two open connected subdomains separated by a
transmission boundary $\Gamma$ that is to say

$$
\Gamma=\bar{D}_{+} \cap \bar{D}_{-}
$$

We denote

$$
D=D_{+} \cup D_{-}=\mathbb{R}^{d} \backslash \Gamma \subset \mathbb{R}^{d}
$$

For a point $x \in \Gamma$ we denote by $\nu(x) \in \mathbb{R}^{d}$ the unit normal to $\Gamma$ at point $x$, pointing to $D_{+}$. In the following, " $f \in C^{p}\left(\bar{D}_{+} ; \mathbb{R}\right) \cap C^{p}\left(\bar{D}_{-} ; \mathbb{R}\right)$ " means that the restriction $f_{+}$of the real valued function $f$ to $D_{+}$(and the restriction $f_{-}$of $f$ to $D_{-}$) coincides on $D_{+}$(resp. $D_{-}$) with a function $\tilde{f}_{+}$of class $C^{p}\left(\mathbb{R}^{d}\right)$ (resp. $\tilde{f}_{-}$). Furthermore $C_{b}^{p}(E)=C_{b}^{p}(E ; \mathbb{R})$ will denote the set of real valued functions on $E$ of class $C^{p}$, bounded with bounded derivatives up to order $p$.

Assume the $a_{i j}$ 's satisfy $\left(a_{ \pm}\right)_{i j} \in C\left(\bar{D}_{ \pm} ; \mathbb{R}\right)$. We may define then the co-normal vector fields $\gamma_{+}(x):=$ $a_{+}(x) \nu(x)$ and $\gamma_{-}(x):=a_{-}(x) \nu(x)$, for $x \in \Gamma$.

We shall consider restricted operators and bilinear forms in the following sense. We define $A_{+}$: $H^{1}\left(D_{+}\right) \rightarrow H^{-1}\left(D_{+}\right)$by

$$
\forall v \in H^{1}\left(D_{+}\right), \quad A_{+} v=\sum_{i, j=1}^{d} D_{i}\left(\left(a_{+}\right)_{i j} D_{j} v\right) .
$$

We define $A_{-}: H^{1}\left(D_{-}\right) \rightarrow H^{-1}\left(D_{-}\right)$in the same manner (note that we do not specify here any domain $\left.\mathcal{D}\left(A_{ \pm}\right)\right)$. Further, we define

$$
\mathcal{E}_{ \pm}(u, v)=\sum_{i, j=1}^{d} \int_{D_{ \pm}}\left(a_{ \pm}\right)_{i j} D_{j} u D_{i} v, \quad \forall u, v \in H^{1}\left(D_{ \pm}\right)
$$

We have, for $u_{ \pm} \in H^{1}\left(D_{ \pm}\right)$with $A_{ \pm} u_{ \pm} \in L^{2}\left(D_{ \pm}\right)$,

$$
\begin{equation*}
\mathcal{E}_{ \pm}\left(u_{ \pm}, v\right)=\int_{D_{ \pm}}\left(-A_{ \pm} u_{ \pm}\right) v, \quad \forall v \in H_{0}^{1}\left(D_{ \pm}\right) \tag{4.1}
\end{equation*}
$$

Imagine now that in (4.1) we wish to take the test function in $H^{1}\left(D_{ \pm}\right)$instead of $H_{0}^{1}\left(D_{ \pm}\right)$. There will still be a link between $A_{ \pm}$and $\mathcal{E}_{ \pm}$, but through Green type identities, involving conormal derivatives and boundary integrals.

We introduce a specific notation for the one-sided conormal derivatives on $\Gamma$ of $u \in L^{2}\left(\mathbb{R}^{d}\right)$ with $u_{ \pm} \in H^{2}\left(D_{ \pm}\right)$. Provided the $\left(a_{ \pm}\right)_{i j}$ are in $C_{b}^{1}\left(\bar{D}_{ \pm} ; \mathbb{R}\right)$ and $\Gamma$ is bounded and Lipschitz we set

$$
\begin{equation*}
\mathcal{B}_{\nu}^{ \pm} u=\nu^{*} \operatorname{Tr}^{ \pm}\left(a_{ \pm} \nabla u_{ \pm}\right)=\sum_{i=1}^{d} \sum_{j=1}^{d} \nu_{i} \operatorname{Tr}^{ \pm}\left(\left(a_{ \pm}\right)_{i j} D_{j} u_{ \pm}\right) \quad \text { on } \quad \Gamma \tag{4.2}
\end{equation*}
$$

where $\operatorname{Tr}^{ \pm}: H^{1}\left(D_{ \pm}\right) \rightarrow H^{1 / 2}(\Gamma)$ stand for the usual trace operators on $\Gamma$.
For $g \in H^{-\frac{1}{2}}(\Gamma)$ and $f \in H^{\frac{1}{2}}(\Gamma)$ we denote by $(g, f)_{\Gamma}$ the action of $g$ on $f$. If both $f, g$ are in $H^{\frac{1}{2}}(\Gamma)$ the quantity $(g, f)_{\Gamma}$ coincides with the surface integral $\int_{\Gamma} g f d s$.

Let us recall the version of the Green identity that is used in the sequel.
Proposition 4.1 (First Green identity, first version; McLean (2000), Lemma 4.1). Assume $\Gamma$ is bounded and $C^{2}$. Let $u \in L^{2}\left(\mathbb{R}^{d}\right)$ with $u_{+} \in H^{2}\left(D_{+}\right)$and $u_{-} \in H^{2}\left(D_{-}\right)$. Assume that the coefficients $\left(a_{ \pm}\right)_{i j}$ are in $C_{b}^{1}\left(\bar{D}_{ \pm} ; \mathbb{R}\right)$. Then

$$
\mathcal{E}_{ \pm}\left(u_{ \pm}, v\right)=\int_{D_{ \pm}}\left(-A_{ \pm} u_{ \pm}\right) v \mp\left(\mathcal{B}_{\nu}^{ \pm} u, \operatorname{Tr}^{ \pm}(v)\right)_{\Gamma}, \quad \forall v \in H^{1}\left(D_{ \pm}\right)
$$

We have the following result.

Theorem 4.2. Assume $\Gamma$ is bounded and $C^{2}$. Assume that Assumption $(\mathbf{E}-\mathbf{B})$ is fulfilled and that for all $1 \leq i, j \leq d$, $a_{i j} \in C_{b}^{1}\left(\bar{D}_{+} ; \mathbb{R}\right) \cap C_{b}^{1}\left(\bar{D}_{-} ; \mathbb{R}\right)$ with $a_{i j}$ possessing a possible discontinuity on $\Gamma$. Then, the Hunt process $\mathbb{M}$ associated to $(\mathcal{E}, \mathcal{D}[\mathcal{E}])$ is a diffusion which possesses the following Skorokhod decomposition: for any $k \in\{1, \ldots, d\}$,

$$
\begin{align*}
X_{t}^{k}= & x_{k}+\int_{0}^{t} \sum_{j=1}^{d} \sigma_{k j}\left(X_{s}\right) d W_{s}^{j}+\int_{0}^{t} \sum_{j=1}^{d} \partial_{j} a_{k j}\left(X_{s}\right) \mathbb{I}_{X_{s} \in D} d s \\
& +\frac{1}{2} \int_{0}^{t} \gamma_{+, k}\left(X_{s}\right) d K_{s}-\frac{1}{2} \int_{0}^{t} \gamma_{-, k}\left(X_{s}\right) d K_{s}, \quad t \geq 0, P^{x}-a . s ., \forall x=\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d} \tag{4.3}
\end{align*}
$$

In the above equality $\sigma: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d \times d}$ denotes the positive square-root of coefficient $2 a$ i.e. the positive matrix real valued coefficient satisfying

$$
\sigma \sigma^{*}(x)=2 a(x), \quad \forall x \in D
$$

(Note that this coefficient exists because $a(x)$ is non-negative definite for all $x \in D$ ). The process $W=\left(W^{1}, \ldots, W^{d}\right)$ is a d-dimensional standard Brownian motion starting from zero and $\left(K_{t}\right)_{t \geq 0}$ denotes the unique PCAF associated to the surface measure $\varsigma(d \xi) \in S$ on $\Gamma$ through the Revuz correspondence. The process $\left(K_{t}\right)$ increases only at times where $X$ lies on $\Gamma$,

$$
\int_{0}^{t} \mathbb{I}_{X_{s} \in \Gamma} d K_{s}=K_{t}, \quad t \geq 0
$$

Proof: We apply the results of Theorem 5.5.5 in Fukushima et al. (1994) in this context for the coordinate functions

$$
p_{k}\left(x_{1}, \ldots, x_{d}\right):=x_{k} \quad(k \in\{1, \ldots, d\})
$$

and follow the ideas of Trutnau (2005) Theorem 5.2. Of course $p_{k} \in \mathcal{D}[\mathcal{E}]_{b, \text { loc }}$ and $p_{k}$ is finely continuous on $\mathbb{R}^{d}$. Let $G$ a relatively compact open set containing $\Gamma$ and a function $f_{k} \in \mathcal{D}[\mathcal{E}]_{b}$ such that $p_{k}=f_{k}$ on $G$. Let $\left\langle M^{\left[f_{k}\right]}\right\rangle$ the square bracket of $M^{\left[f_{k}\right]}$. Then, an easy computation from (3.3) shows that the energy measure of $M^{\left[f_{k}\right]}$ (the Revuz measure of $\left\langle M^{\left[f_{k}\right]}\right\rangle$ ) is

$$
\begin{equation*}
\mu_{\left\langle f_{k}\right\rangle}(d y)=\mu_{\left\langle M^{\left.\left[f_{k}\right]\right\rangle}\right.}(d y)=\left\langle 2 a(y) \nabla f_{k}(y), \nabla f_{k}(y)\right\rangle \ell(d y) \tag{4.4}
\end{equation*}
$$

and we know that $\mu_{\left\langle f_{k}\right\rangle}=\mu_{\left\langle p_{k}\right\rangle}$ on $G$. It is easy to show that $\mathbb{I}_{G} \cdot \mu_{\left\langle p_{k}\right\rangle}$ is a finite Radon measure belonging to $S_{00}$ and that $\mu_{\left\langle p_{k}\right\rangle}$ is a smooth measure. Then, an easy computation from (3.2) shows that

$$
\left\langle M^{\left[f_{k}\right]}\right\rangle_{t}=\int_{0}^{t}\left\langle 2 a\left(X_{s}\right) \nabla f_{k}\left(X_{s}\right), \nabla f_{k}\left(X_{s}\right)\right\rangle d s, \quad k \in\{1, \ldots, d\}
$$

and by the well-known results on stochastic representation of martingales, there exists a dimensional Brownian motion $W=\left(W^{1}, \ldots, W^{d}\right)$ such that

$$
M_{t}^{\left[f_{k}\right]}=\int_{0}^{t}\left[\sigma\left(X_{s}\right) \nabla f_{k}\left(X_{s}\right)\right]^{*} d W_{s}, \quad P^{x}-\text { a.s. } \forall x \in \mathbb{R}^{d}, \quad k \in\{1, \ldots, d\}
$$

(see for e.g. Revuz and Yor (1999) Chapter V Theorem 3.9 and the remark following its proof).

Moreover, for any $v \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, using the Green Identities of Proposition 4.1 and taking into account that $v$ is of compact support, we have:

$$
\begin{aligned}
\mathcal{E}\left(f_{k}, v\right)= & \mathcal{E}_{+}\left(f_{k,+}, v\right)+\mathcal{E}_{-}\left(f_{k,-}, v\right) \\
= & \int_{D_{+}}\left(-A_{+} f_{k,+}\right) v-\left(\mathcal{B}_{\nu}^{+} f_{k,+}, \operatorname{Tr}^{+}(v)\right)_{\Gamma}+\int_{D_{-}}\left(-A_{-} f_{k,-}\right) v+\left(\mathcal{B}_{\nu}^{-} f_{k,-}, \operatorname{Tr}^{-}(v)\right)_{\Gamma} \\
= & -\int_{D} \sum_{i, j=1}^{d} D_{i}\left(a_{i j}(y) D_{j} f_{k}(y)\right) v(y) \mathbb{I}_{y \in D} \ell(d y) \\
& \quad-\int_{\Gamma} \nu^{*}\left[\operatorname{Tr}^{+}\left(a_{+} \nabla f_{k,+}\right)-\operatorname{Tr}^{-}\left(a_{-} \nabla f_{k,-}\right)\right] \gamma(v) d \varsigma \\
= & -\int_{D} \sum_{j=1}^{d} \partial_{j} a_{k j}(y) v(y) \mathbb{I}_{y \in D} \ell(d y)-\int_{\Gamma}\left[\operatorname{Tr}^{+}\left(\left(a_{+} \nu\right)_{k}\right)-\operatorname{Tr}^{-}\left(\left(a_{-} \nu\right)_{k}\right)\right] v d \varsigma \\
= & -\int_{D} \sum_{j=1}^{d} \partial_{j} a_{k j}(y) v(y) \mathbb{I}_{y \in D} \ell(d y)-\int_{\Gamma}\left[\left(\tilde{a}_{+} \nu\right)_{k}-\left(\tilde{a}_{-} \nu\right)_{k}\right] v d \varsigma \\
= & \left(\varrho_{k}^{+}, v\right)-\left(\varrho_{k}^{-}, v\right)
\end{aligned}
$$

with

$$
\varrho_{k}^{ \pm}(d y):=-\sum_{j=1}^{d}\left[\partial_{j} a_{k j}(y)\right]^{ \pm} \mathbb{I}_{y \in D} \ell(d y)+\left[\left(\gamma_{-}\right)_{k}-\left(\gamma_{+}\right)_{k}\right]^{ \pm}(y) \mathbb{I}_{y \in \Gamma} \varsigma(d y)
$$

(here, the notation $[a]^{+}$(resp. $[a]^{-}$) stands for the positive (resp. negative) part of some real number $a$ ).

Let us now proceed to show that the measures $\mathbb{I}_{G} \cdot \varrho_{k}^{ \pm}$belong to $S_{00}$.
Note that $\left\|\partial_{j} a_{k j} \mathbb{I}_{D \cap G}\right\|_{\infty}<\infty$ and from the definition of $S_{00}$ and the Revuz correpondence (3.2), it is not difficult to prove that the measures $-\left[\partial_{j} a_{k j}\right]_{ \pm}(y) \mathbb{I}_{y \in D} \ell(d y)$ are smooth with their corresponding additive functional writing as $\left(-\int_{0}^{t}\left[\partial_{j} a_{k j}\right]_{ \pm}\left(X_{s}\right) \mathbb{I}_{X_{s} \in D} d s\right)_{t \geq 0}$.

We now turn to the surface measures $\zeta_{k}^{ \pm}(d y):=\left[\left(\gamma_{-}\right)_{k}-\left(\gamma_{+}\right)_{k}\right]^{ \pm}(y) \mathbb{I}_{y \in \Gamma} \varsigma(d y)$. It is well-known (see e.g. Evans and Gariepy (2015) p. 134 3. ( $\star \star \star)$ ) ( $(\star \star \star \star)$ ) that there exists a universal constant $C_{0}>0$, depending only on the Lipschitz domain $D_{+}$, such that for all $h \in C^{1}\left(\bar{D}_{+}\right)$,

$$
\int_{\Gamma}|h(y)| \varsigma(d y) \leq C_{0} \int_{D_{+}}(|\nabla h(x)|+|h(x)|) \ell(d x) .
$$

Thus, for all $h \in \mathcal{D}[\mathcal{E}] \cap C_{c}\left(\mathbb{R}^{d}\right)$, we have

$$
\begin{aligned}
\int_{\Gamma}|h(y)| \varsigma(d y) & \leq C_{0} \int_{D_{+}}(|\nabla h(x)|+|h(x)|) \ell(d x) \\
& \leq C_{0}\left(2 \ell\left(D_{+}\right)\right)^{1 / 2}\left(\int_{\mathbb{R}^{d}}\left(|\nabla h(x)|^{2}+|h(x)|^{2}\right) \ell(d x)\right)^{1 / 2} \\
& \leq C_{0} \sqrt{\frac{\left(2 \ell\left(D_{+}\right)\right)}{\lambda}}(\mathcal{E}(h, h)+(h, h))^{1 / 2}
\end{aligned}
$$

so that the surface measure $\varsigma(d y)$ belongs to $S_{0}$. Since

$$
\forall y \in \Gamma, \quad\left[\left|\left(\gamma_{-}\right)_{k}-\left(\gamma_{+}\right)_{k}\right|(y)\right]^{ \pm} \leq 2\left|\tilde{a}_{ \pm}(y) \nu(y)\right| \leq 2 \Lambda,
$$

the surface measures $\zeta_{k}^{ \pm}(d y):=\left[\left(\gamma_{-}\right)_{k}-\left(\gamma_{+}\right)_{k}\right]^{ \pm}(y) \mathbb{I}_{y \in \Gamma} \varsigma(d y)$ belong also to $S_{0}$.

Note that from Aronson's estimates (2.13) we retrieve the following estimations

$$
r_{1}(x, y) \leq C|x-y|^{-(d-2)} \quad \text { if } d>2 ; r_{1}(x, y) \leq C(\ln (1 /|x-y|) \vee 1) \text { if } d=2
$$

Then, using the same arguments as in Fukushima et al. (1994) (Example 5.2 .2 p.255), we can assert that the measures $\zeta_{k}^{ \pm}(d y)$ belong to $S_{00}$. Moreover, let $\left(K_{t}\right)_{t \geq 0}$ denote the PCAF associated to $\varsigma(d y)$; in regard of the results stated in the original article of D. Revuz (cf. Revuz (1970) p.507) we may assert that $\left(\int_{0}^{t}\left[\left(\gamma_{-}\right)_{k}-\left(\gamma_{+}\right)_{k}\right]^{ \pm}\left(X_{s}\right) \mathbb{I}_{X_{s} \in \Gamma} d K_{s}\right)_{t \geq 0}$ is the PCAF associated to $\zeta_{k}^{ \pm}(d y)$ via the Revuz correspondence.

By application of Theorem 5.5.5 in Fukushima et al. (1994) and since all the necessary hypothesis are fulfilled, we get the decomposition (4.3) on the set $\left\{t \geq 0: t \leq \tau_{G_{q}}\right\}$ where $G_{q}:=\{x \in$ $\left.\mathbb{R}^{d}:|x|<q\right\}$. The identification of the process for all times follows by letting $q$ tend to infinity.

Let $u_{0} \in \mathcal{D}(A)$. From the Hille-Yosida theorem (Brezis (1983) Theorems VII. 4 and VII.5) we can prove that there exists a unique function

$$
u \in C^{1}\left([0, T] ; L^{2}\left(\mathbb{R}^{d}\right)\right) \cap C([0, T] ; \mathcal{D}(A))
$$

satisfying

$$
\begin{equation*}
\frac{\mathrm{d} u}{\mathrm{~d} t}=A u, \quad u(0)=u_{0} \tag{4.5}
\end{equation*}
$$

where the first equality in (4.5) has to be understood in the weak sense.
Under the hypothesis of Theorem 4.2, we deduce the following Corollary.
Corollary 4.3. Let $0<T<\infty$. Under the conditions of Theorem 4.2, for any $u_{0} \in \mathcal{D}(A)$, we have

$$
\begin{equation*}
E^{x}\left[u_{0}\left(X_{t}\right)\right]=u(t, x), \quad \forall t \in[0, T], \forall x \in \mathbb{R}^{d} \tag{4.6}
\end{equation*}
$$

where $X$ is the diffusion considered in Theorem 4.2 and $u$ is the solution of (4.5).
In particular, the following transmission condition

$$
\begin{equation*}
\left\langle a_{+} \nabla_{x} u_{+}(t, y)-a_{-} \nabla_{x} u_{-}(t, y), \nu(y)\right\rangle=0, \text { for a.e. }(t, y) \in(0, T] \times \Gamma \quad(\star) \tag{4.7}
\end{equation*}
$$

is satisfied.
Proof: In view of (3.1) and since $\frac{d}{d t} T_{t} u_{0}=A T_{t} u_{0}$ (Pazy (1983) Thm 2.4-c)) the function $(t, x) \mapsto$ $E^{x}\left[u_{0}\left(X_{t}\right)\right]$ is solution of (4.5). We refer to the proofs of Proposition 3.14 and Theorem 3.1 in Étoré and Martinez (2021) for the verification of the other assertions.

In the light of (4.6) and in order to compute an approximate value of $u(t, x)$, one could think of producing a Monte Carlo method. Our article Étoré and Martinez (2021) is an attempt to tackle this issue.
4.2. The diagonal case: link with the results of Bossy $छ$ al. Bossy et al. (2010). We wish to compare the result of our Theorem 4.2 with the ones in Bossy et al. (2010). For this purpose we restrict once more the assumption on the diffusion coefficient $a(x)$ given in the setting of the preceding subsection. Namely we assume

$$
a(x)=I_{d} \varepsilon(x)
$$

with $\varepsilon(x):=\left[\varepsilon_{+} \mathbf{1}_{\bar{D}_{+}}+\varepsilon_{-} \mathbf{1}_{D_{-}}\right](x), I_{d}$ the identity matrix and $\varepsilon_{+} \neq \varepsilon_{-} \in \mathbb{R}_{+}^{*}$. We are therefore considering the case of a diagonal diffusion matrix. We will also assume in the forthcoming theorems that $\Gamma$ is bounded and closed, and will consider that $D_{+}$is the interior domain delimited by $\Gamma$.

We start by summing up the notions and results in Bossy et al. (2010) that we need in order to do the comparison.

In the sequel we denote $(C, \mathcal{C})$ the usual canonical space, i.e. $C=C\left([0, \infty) ; \mathbb{R}^{d}\right)$ and $\mathcal{C}=\mathcal{B}(C)$ (see Pb 2.4 .2 in Karatzas and $\operatorname{Shreve}(1991)$ for details). We also denote $\left(\mathcal{C}_{t}\right)_{t \geq 0}$ the usual canonical filtration (see Eq. (5.3.19) in Karatzas and Shreve (1991)).

We shall denote by $\omega$ the canonical process defined on $(C, \mathcal{C})$. Note that $\omega=\left(\omega_{t}\right)_{t \geq 0}$ is $\left(\mathcal{C}_{t}\right)$ adapted.

Let us define the transmission operator $\mathcal{L}$ acting on functions $\varphi \in C\left(\mathbb{R}^{d}\right) \cap C^{2}\left(\bar{D}_{+}\right) \cap C^{2}\left(\bar{D}_{-}\right)$by

$$
\begin{align*}
\mathcal{L} \varphi(x) & =\sum_{i, j=1}^{d} \partial_{i}\left(\left[I_{d} \varepsilon\right]_{i j} \partial_{j} \varphi\right)(x)=\sum_{i=1}^{d} \varepsilon(x) \partial_{i i}^{2} \varphi(x), \quad \forall x \in \mathbb{R}^{d} \backslash \Gamma, \\
\mathcal{L} \varphi(x) & =\delta(x), \quad \forall x \in \Gamma \tag{4.8}
\end{align*}
$$

where for any $x \in \Gamma, \delta(x) \in \mathbb{R}$ is an arbitrary value of no importance in our computations.
Definition 4.4 (Bossy et al. (2010)). i ) A family of probability measures $\left(\mathbb{P}^{x}\right)_{x \in \mathbb{R}^{d}}$ on $(C, \mathcal{C})$ solves the martingale problem for the operator $\mathcal{L}$ if, for all $x \in \mathbb{R}^{d}$ one has $\mathbb{P}^{x}(\omega(0)=x)=1$ and, for all $\varphi$ satisfying

$$
\begin{gather*}
\varphi \in C_{b}\left(\mathbb{R}^{d}\right) \cap C_{b}^{2}\left(\bar{D}_{+}\right) \cap C_{b}^{2}\left(\bar{D}_{-}\right)  \tag{4.9}\\
\left\langle\varepsilon_{+} \nabla_{x} \varphi_{+}(y)-\varepsilon_{-} \nabla_{x} \varphi_{-}(y), \nu(y)\right\rangle=0, \quad \forall y \in \Gamma, \quad(\star) \tag{4.10}
\end{gather*}
$$

one has that the $t$-indexed process defined by

$$
\varphi\left(\omega_{t}\right)-\varphi\left(\omega_{0}\right)-\int_{0}^{t} \mathcal{L} \varphi\left(\omega_{s}\right) d s, \quad \forall t \geq 0
$$

is a $\left(\mathcal{C}_{t}\right)$-martingale under $\mathbb{P}^{x}$.
ii) The martingale problem is said to be well-posed if there exists a unique family of probability measures $\left(\mathbb{P}^{x}\right)_{x \in \mathbb{R}^{d}}$ which solves the martingale problem for the operator $\mathcal{L}$.

Theorem 4.5 (Theorems 2.4, 2.10 and 2.12 in Bossy et al. (2010)). For all $x \in \mathbb{R}^{d}$ consider the SDE

$$
\left\{\begin{array}{l}
X_{t}=x+\int_{0}^{t} \sqrt{2 \varepsilon\left(X_{u}\right)} d B_{u}+\frac{\varepsilon_{+}-\varepsilon_{-}}{2 \varepsilon_{-}} \int_{0}^{t} \nu\left(X_{u}\right) d L_{u}^{0}(Y)  \tag{4.11}\\
Y_{t}=\rho\left(X_{t}\right)
\end{array}\right.
$$

where $B$ is a d-dimensional Brownian motion, $\rho(z)$ is the distance from $z \in \mathbb{R}^{d}$ to the boundary $\Gamma$ and $L_{t}^{0}(Y)$ stands for the (right) local time at point zero of the local martingale $Y$.

If $\Gamma$ is of class $C^{3}$ and compact we have:
i) There exists a weak solution to (4.11). For any $x \in \mathbb{R}^{d}$ consider the law $\mathbb{P}^{x}$ of this weak solution on $(C, \mathcal{C})$. The family $\left(\mathbb{P}^{x}\right)$ is a solution to the martingale problem for $\mathcal{L}$.
ii) Conversely let $\left(\mathbb{P}^{x}\right)$ be a solution to the martingale problem for $\mathcal{L}$, and let $x \in \mathbb{R}^{d}$. There exists a $\left(\mathcal{C}_{t}\right)$-Brownian motion $B$ under $\mathbb{P}^{x}$, such that $\omega$ is a weak solution to (4.11) driven by $B$, under $\mathbb{P}^{x}$.

Remark 4.6. Note a change of sign in the weight in front of the local time term in (4.11), compared to Bossy et al. (2010). This is because for us the exterior normal to the interior domain $D_{+}$is $-\nu$.

Remark 4.7. Note that in Bossy et al. (2010) the authors work with right local time at point zero of $Y=\rho(X)$. Working with the symmetric local time instead, as it is often the case in the study of asymetric diffusions (e.g. Lejay (2006)) would lead to different coefficients in front of the local time term.

Theorem 4.8 (Theorem 2.14 in Bossy et al. (2010)). Assume $\Gamma$ is of class $C^{3}$ and compact. Then the martingale problem for $\mathcal{L}$ is well posed and in particular there is a unique weak solution to (4.11) in the sense of probability law.

We then have the following main result.

Theorem 4.9. Assume $\Gamma$ is of class $C^{2}$ and compact. Consider the Hunt process $X(\mathbb{M})$ in Theorem 4.2. For any starting point $x \in \mathbb{R}^{d}$ consider the law $\mathbb{P}^{x}$ of $X$ on $(C, \mathcal{C})$. The family $\left(\mathbb{P}^{x}\right)$ is a solution to the martingale problem for $\mathcal{L}$.

Proof: Let $\mathbb{M}=\left(\Omega,\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathcal{F},\left(X_{t}\right)_{t \geq 0},\left(P^{x}\right)_{x \in \mathbb{R}^{d}}\right)$ the Hunt process considered in Theorem 4.2. Let $x \in \mathbb{R}^{d}$ and let $\varphi$ satisfying (4.9) and (4.10). We aim at applying Theorem 5.5.5. in Fukushima et al. (1994) with $\varphi$ in order to check that

$$
\varphi\left(X_{t}\right)-\varphi(x)-\int_{0}^{t} \mathcal{L} \varphi\left(X_{s}\right) d s
$$

is a martingale under $P^{x}$ w.r.t $\left(\mathcal{F}_{t}\right)$.
Indeed, proceeding as in Karatzas and Shreve (1991) p314, this implies that the law $\mathbb{P}^{x}$ induced on $(C, \mathcal{C})$ by $X$ (under $\left.P^{x}\right)$ is such that if we consider the family $\left(\mathbb{P}^{x}\right)$ this is a solution to the martingale problem for $\mathcal{L}$.

The fact that $\varphi$ is finely continuous is clear. In order to check that $\varphi \in \mathcal{D}[\mathcal{E}]_{b, \text { loc }}$ it suffices to notice that for any compact $G \subset \mathbb{R}^{d}$ the functions $\left(\varphi \mathbf{1}_{G}\right)_{+}$and $\left(\varphi \mathbf{1}_{G}\right)_{+}$are respectively in $H^{1}(D+)$ and $H^{1}\left(D_{-}\right)$. As $\varphi \mathbf{1}_{G}$ is in $L^{2}\left(\mathbb{R}^{d}\right)$ and continuous across $\Gamma$, Exercise 4.5 in McLean (2000) implies that $\varphi \mathbf{1}_{G} \in H^{1}\left(\mathbb{R}^{d}\right)=\mathcal{D}[\mathcal{E}]$.

We now check that $\mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle} \in S_{00}$ for any relatively compact set $G$. We first check that $\mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}$ is a positive and finite Radon measure. Indeed one has, as in Eq. (3.3),

$$
\mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}\left(\mathbb{R}^{d}\right)=\int_{G} \mu_{\langle\varphi\rangle}(d x)=2 \mathcal{E}\left(\mathbf{1}_{G} \varphi, \varphi\right)-\mathcal{E}\left(\mathbf{1}_{G} \varphi^{2}, 1\right)=2 \mathcal{E}\left(\mathbf{1}_{G} \varphi, \varphi\right)
$$

and $0 \leq \mathcal{E}\left(\mathbf{1}_{G} \varphi, \varphi\right)<\infty$. We now check that $\left\|R_{1} \mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}\right\|_{\infty}<\infty$. For any $x \in \mathbb{R}^{d}$ we have, with a constants $m, M^{\prime}$ depending on $\lambda, \Lambda$,

$$
\begin{align*}
\left|R_{1} \mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}(x)\right| & \leq \int_{0}^{\infty} e^{-t}\left(\int_{\mathbb{R}^{d}} p(t, x, y)\left|\mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}\right|(d y)\right) d t \\
& \leq M \int_{0}^{\infty} d t e^{-t} \int_{\mathbb{R}^{d}} \frac{1}{t^{d / 2}} e^{-\frac{|y-x|^{2}}{M t}}\left|\mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}\right|(d y) \\
& =M \int_{0}^{\infty} d t e^{-t} \int_{G} \frac{1}{t^{d / 2}} e^{-\frac{|y-x|^{2}}{M t}}|\langle 2 a(y) \nabla \varphi(y), \nabla \varphi(y)\rangle| d y  \tag{4.12}\\
& =M^{\prime} \int_{0}^{\infty} d t e^{-t} E^{x}\left[\left|\left\langle 2 a\left(m W_{t}\right) \nabla \varphi\left(m W_{t}\right), \nabla \varphi\left(m W_{t}\right)\right\rangle\right| \mathbf{1}_{W_{t} \in G}\right] \\
& \leq 2 M^{\prime} \Lambda_{z \in G}|\nabla u(m z)|^{2} \int_{0}^{\infty} d t e^{-t} d t<\infty
\end{align*}
$$

Here we have used (2.13) at the second line. At the third line we have used (3.3) and computations similar to the ones leading to (4.4). At the last line we have used (E-B). As the bound in (4.12) does not depend on $x$ we have proven $\left\|R_{1} \mathbf{1}_{G} \cdot \mu_{\langle\varphi\rangle}\right\|_{\infty}<\infty$.

Let $v \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$. We have, using in particular Proposition 4.1, the smoothness of $\varphi, \operatorname{Tr}^{+}(v)=$ $\operatorname{Tr}(v+)=\operatorname{Tr}^{-}(v)=\operatorname{Tr}(v-)=\operatorname{Tr}(v)$ and (4.10),

$$
\begin{aligned}
\mathcal{E}(\varphi, v) & =\mathcal{E}_{+}(\varphi, v)+\mathcal{E}_{-}(\varphi, v) \\
& =-\int_{D_{+}}\left(A_{+} \varphi_{+}\right) v-\int_{D_{-}}\left(A_{-} \varphi_{-}\right) v+\left(\mathcal{B}_{\nu}^{-} u-\mathcal{B}_{\nu}^{+} u, \operatorname{Tr}(v)\right)_{\Gamma} \\
& =-\int_{D} \mathcal{L} \varphi v+\int_{\Gamma}\left\langle\operatorname{Tr}\left(\varepsilon_{+} \nabla_{x} \varphi_{+}\right)-\operatorname{Tr}\left(\varepsilon_{-} \nabla_{x} \varphi_{-}\right), \nu\right\rangle v d \varsigma \\
& =-\int_{D} \mathcal{L} \varphi v .
\end{aligned}
$$

The function $\mathcal{L} \varphi$ is piecewise continuous and bounded and proceeding as above one may check that the positive and negative parts $(-\mathcal{L} \varphi)^{ \pm}$satisfy $\mathbf{1}_{G} \cdot\left((-\mathcal{L} \varphi)^{ \pm} d x\right) \in S_{00}$ for any relatively compact set $G$.

We denote $A^{ \pm}$the PCAF's related to Revuz measures $(-\mathcal{L} \varphi)^{ \pm} d x$. We set

$$
N^{[\varphi]}=-A^{+}+A^{-}
$$

and notice that following Example 5.1.1 in Fukushima et al. (1994) we have $N_{t}^{[\varphi]}=-\int_{0}^{t}(-\mathcal{L} \varphi)^{+}\left(X_{s}\right) d s+$ $\int_{0}^{t}(-\mathcal{L} \varphi)^{-}\left(X_{s}\right) d s=\int_{0}^{t} \mathcal{L} \varphi\left(X_{s}\right) d s$ (note that we use the fact that $\int_{0}^{t} \mathbf{1}_{X_{s} \in \Gamma} d s=0 \mathbb{P}^{x}$-a.s., therefore the arbitrary value of $\mathcal{L} \varphi(z)$ for $z \in \Gamma$ causes no issue).

We now apply Theorem 5.5.5. in Fukushima et al. (1994). We have

$$
\varphi\left(X_{t}\right)-\varphi(x)=M_{t}^{[\varphi]}+\int_{0}^{t} \mathcal{L} \varphi\left(X_{s}\right) d s
$$

with $M^{[\varphi]}$ which a martingale (as for example the $M^{\left[f_{k}\right]}$ 's are martingales in the proof of Theorem 4.2). The proof is completed.

In view of Theorems 4.5, 4.8 and 4.9 we immediately get the following corollary.
Corollary 4.10. Assume $\Gamma$ is of class $C^{3}$ and compact.
The solutions of (4.3) and (4.11), with starting point $x \in \mathbb{R}^{d}$, have the same distribution on (C, $\mathcal{C}$ ).

In particular there is uniqueness in the sense of probability law of the weak solutions of (4.3).
Going a bit further in the analysis we may do an identification in the strong sense of the terms appearing in (4.3) and (4.11).
Corollary 4.11. Assume $\Gamma$ is of class $C^{3}$ and compact.
Consider the Hunt process $\mathbb{M}=\left(\Omega,\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathcal{F},\left(X_{t}\right)_{t \geq 0},\left(P^{x}\right)_{x \in \mathbb{R}^{d}}\right)$ in Theorem 4.2, which is such that $X$ solves (4.3) under $P^{x}$.

Then one also has

$$
\left\{\begin{array}{l}
X_{t}=x+\int_{0}^{t} \sqrt{2 \varepsilon\left(X_{u}\right)} d B_{u}(X)+\frac{\varepsilon_{+}-\varepsilon_{-}}{2 \varepsilon_{-}} \int_{0}^{t} \nu\left(X_{u}\right) d L_{u}^{0}(Y) \\
Y_{t}=\rho\left(X_{t}^{0}\right.
\end{array}\right.
$$

with $B(X)$ a Brownian motion under $P^{x}$ (that is $X$ is also a weak solution to (4.11)).
Moreover one has $B(X)=W$ (with $W$ the Brownian motion driving (4.3)) and

$$
\begin{equation*}
K_{t}=\frac{1}{\varepsilon_{-}} L_{t}^{0}(\rho(X)), \quad \forall t \geq 0 . \tag{4.13}
\end{equation*}
$$

Proof: Pick $x \in \mathbb{R}^{d}$. Let us rewrite (4.3) in the matrix form in the case of interest. One has

$$
X_{t}=x+\int_{0}^{t} \sqrt{2 \varepsilon}\left(X_{u}\right) d W_{u}+\frac{1}{2} \int_{0}^{t}\left(\gamma_{+}\left(X_{u}\right)-\gamma_{-}\left(X_{u}\right)\right) d K_{u}
$$

under $P^{x}$. We set now

$$
G(X):=X_{t}-x=\int_{0}^{t} \sqrt{2 \varepsilon}\left(X_{u}\right) d W_{u}+\frac{1}{2} \int_{0}^{t}\left(\gamma_{+}\left(X_{u}\right)-\gamma_{-}\left(X_{u}\right)\right) d K_{u}
$$

and notice that $G(\cdot)$ is obviously measurable and that

$$
\Phi(X):=G(X)-X_{t}-x=0 \quad P^{x}-\text { a.s. }
$$

Noting that $\Phi(\cdot)$ is measurable and remembering the definition of $\mathbb{P}^{x}$ as the law of $X$ on $(C, \mathcal{C})$ under $P^{x}$ it is obvious that one has $P^{x}(\Phi(X)=0)=1=\mathbb{P}^{x}(\Phi(\omega)=0)$.

Therefore under $\mathbb{P}^{x}$ one has a.s.

$$
G(\omega)=\omega_{t}-x
$$

but from Theorem 4.5 Point ii) one has

$$
\begin{equation*}
G(\omega)=\int_{0}^{t} \sqrt{2 \varepsilon\left(\omega_{u}\right)} d B_{u}(\omega)+\frac{\varepsilon_{+}-\varepsilon_{-}}{2 \varepsilon_{-}} \int_{0}^{t} \nu\left(\omega_{u}\right) d L_{u}^{0}(\rho(\omega)) \tag{4.14}
\end{equation*}
$$

where $B(\omega)$ is a $\left(\mathcal{C}_{t}\right)$-Brownian motion under $\mathbb{P}^{x}$. In the notation we have stressed that $B(\omega)$ is constructed from the paths of $\omega$ (through a measurable mapping).

Applying $G(\cdot)$ viewed as in (4.14) to $X$ we get the first part of the corollary, that is that $X$ solves (4.11) (4.11) driven by $B(X)$ (to check that $B(X)$ is a B.m. under $P^{x}$ we have to check that the increments of $B(X)$ are independent and that $B_{t}(X)-B_{s}(X)$ for any $s<t$ and distributed as a $\mathcal{N}_{d}\left(0,(t-s) I_{d}\right)$, by identifying the law of $X$ under $P^{x}$ with the one of $\omega$ under $\left.\mathbb{P}^{x}\right)$.

Using the uniqueness of the decomposition of a semimartingale and recalling that $\gamma_{ \pm}=a_{ \pm} \nu$ we have

$$
\begin{equation*}
\int_{0}^{t} \sqrt{2 \varepsilon}\left(X_{u}\right) d W_{u}=\int_{0}^{t} \sqrt{2 \varepsilon\left(X_{u}\right)} d B_{u}(X) \quad \forall t \geq 0 \tag{4.15}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{0}^{t} \nu\left(X_{s}\right) d K_{s}=\int_{0}^{t} \nu\left(X_{s}\right) \frac{1}{\varepsilon_{-}} d L_{s}^{0}(\rho(X)) \quad \forall t \geq 0 \tag{4.16}
\end{equation*}
$$

From (4.15) and (1.1) one has

$$
\left.0 \leq\langle W-B(X)\rangle_{t} \leq \frac{1}{2 \lambda} \int_{0}^{t} 2 \varepsilon\left(X_{u}\right) d\langle W-B(X)\rangle_{u} \leq \frac{1}{2 \lambda}\left\langle\int_{0} \sqrt{2 \varepsilon}\left(X_{u}\right) d(W-B(X))_{u}\right)\right\rangle_{t}=0
$$

where the bracket has to be understood in the multidimensional sense (matrix of brackets). Therefore $W \equiv B(X)$ using Proposition IV.1.12 in Revuz and Yor (1999).

It remains to use the componentwise meaning of (4.16) in order to check (4.13). For the use of exposure we assume for a while that $d=2$, with $\nu=\left(\nu_{1}, \nu_{2}\right)^{T}$. Using the fact that if $\nu_{1}(y)=0$ then
$\nu_{2}(y) \neq 0, y \in \Gamma$, we get for any $t \geq 0$,

$$
\begin{aligned}
K_{t} & =\int_{0}^{t}\left(\mathbf{1}_{\nu_{1}\left(X_{s}\right) \neq 0} \frac{\nu_{1}\left(X_{s}\right)}{\nu_{1}\left(X_{s}\right)}+\mathbf{1}_{\nu_{1}\left(X_{s}\right)=0} \frac{\nu_{2}\left(X_{s}\right)}{\nu_{2}\left(X_{s}\right)}\right) d K_{s} \\
& =\int_{0}^{t}\left(\mathbf{1}_{\nu_{1}\left(X_{s}\right) \neq 0} \frac{\nu_{1}\left(X_{s}\right)}{\nu_{1}\left(X_{s}\right)}+\mathbf{1}_{\nu_{1}\left(X_{s}\right)=0} \frac{\nu_{2}\left(X_{s}\right)}{\nu_{2}\left(X_{s}\right)}\right) \frac{1}{\varepsilon_{-}} d L_{s}^{0}(\rho(X)) \\
& =\int_{0}^{t}\left(\mathbf{1}_{\nu_{1}\left(X_{s}\right) \neq 0}+\mathbf{1}_{\nu_{1}\left(X_{s}\right)=0}\right) \frac{1}{\varepsilon_{-}} d L_{s}^{0}(\rho(X)) \\
& =\frac{1}{\varepsilon_{-}} L_{t}^{0}(\rho(X))
\end{aligned}
$$

We claim that the above reasoning can easily be extended to $d>2$. The proof is completed.

## Acknowledgements

The authors would like to thank the anonymous referee for his valuable remarks about preliminary versions of this paper.

## References

Aronson, D. G. Bounds for the fundamental solution of a parabolic equation. Bull. Amer. Math. Soc., 73, 890-896 (1967). doi:10.1090/S0002-9904-1967-11830-5. URL https://doi.org/10. 1090/S0002-9904-1967-11830-5.
Bass, R. F. Diffusions and elliptic operators. Probability and its Applications (New York). SpringerVerlag, New York (1998). ISBN 0-387-98315-5.
Bossy, M., Champagnat, N., Maire, S., and Talay, D. Probabilistic interpretation and random walk on spheres algorithms for the Poisson-Boltzmann equation in molecular dynamics. M2AN Math. Model. Numer. Anal., 44 (5), 997-1048 (2010). doi:10.1051/m2an/2010050. URL https: //doi.org/10.1051/m2an/2010050.
Brezis, H. Analyse fonctionnelle. Collection Mathématiques Appliquées pour la Maîtrise. [Collection of Applied Mathematics for the Master's Degree]. Masson, Paris (1983). ISBN 2-225-77198-7. Théorie et applications. [Theory and applications].
Chung, K. L. and Walsh, J. B. To reverse a Markov process. Acta Math., 123, 225-251 (1969). doi:10.1007/BF02392389. URL https://doi.org/10.1007/BF02392389.
Étoré, P. and Martinez, M. A transformed stochastic Euler scheme for multidimensional transmission PDE. J. Comput. Appl. Math., 394, Paper No. 113551, 28 (2021). doi:10.1016/j.cam. 2021. 113551. URL https://doi.org/10.1016/j.cam.2021.113551.

Evans, L. C. and Gariepy, R. F. Measure theory and fine properties of functions. Textbooks in Mathematics. CRC Press, Boca Raton, FL, revised edition (2015). ISBN 978-1-4822-4238-6.
Fukushima, M., Ōshima, Y., and Takeda, M. Dirichlet forms and symmetric Markov processes, volume 19 of De Gruyter Studies in Mathematics. Walter de Gruyter \& Co., Berlin (1994). ISBN 3-11-011626-X. doi:10.1515/9783110889741. URL https://doi.org/10.1515/9783110889741.
Haussmann, U. G. and Pardoux, E. Time reversal of diffusions. Ann. Probab., 14 (4), 1188-1205 (1986). URL http://links.jstor.org/sici?sici=0091-1798(198610)14:4<1188:TROD>2.0. CO;2-Z\&origin=MSN.
Karatzas, I. and Shreve, S. E. Brownian motion and stochastic calculus, volume 113 of Graduate Texts in Mathematics. Springer-Verlag, New York, second edition (1991). ISBN 0-387-97655-8. doi:10.1007/978-1-4612-0949-2. URL https://doi.org/10.1007/978-1-4612-0949-2.

Lejay, A. On the constructions of the skew Brownian motion. Probab. Surv., 3, 413-466 (2006). doi:10.1214/154957807000000013. URL https://doi.org/10.1214/154957807000000013.
Lyons, T. J. and Zheng, W. A. Diffusion processes with nonsmooth diffusion coefficients and their density functions. Proc. Roy. Soc. Edinburgh Sect. A, 115 (3-4), 231-242 (1990). doi: 10.1017/S0308210500020618. URL https://doi.org/10.1017/S0308210500020618.

McLean, W. Strongly elliptic systems and boundary integral equations. Cambridge University Press, Cambridge (2000). ISBN 0-521-66332-6; 0-521-66375-X.
Pazy, A. Semigroups of linear operators and applications to partial differential equations, volume 44 of Applied Mathematical Sciences. Springer-Verlag, New York (1983). ISBN 0-387-90845-5. doi: 10.1007/978-1-4612-5561-1. URL https://doi.org/10.1007/978-1-4612-5561-1.

Revuz, D. Mesures associées aux fonctionnelles additives de Markov. I. Trans. Amer. Math. Soc., 148, 501-531 (1970). doi:10.2307/1995386. URL https://doi.org/10.2307/1995386.
Revuz, D. and Yor, M. Continuous martingales and Brownian motion, volume 293 of Grundlehren der mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. Springer-Verlag, Berlin, third edition (1999). ISBN 3-540-64325-7. doi:10.1007/ 978-3-662-06400-9. URL https://doi.org/10.1007/978-3-662-06400-9.
Rozkosz, A. Stochastic representation of diffusions corresponding to divergence form operators. Stochastic Process. Appl., 63 (1), 11-33 (1996a). doi:10.1016/0304-4149(96)00059-2. URL https: //doi.org/10.1016/0304-4149(96)00059-2.
Rozkosz, A. Weak convergence of diffusions corresponding to divergence form operators. Stochastics Stochastics Rep., 57 (1-2), 129-157 (1996b). doi:10.1080/17442509608834055. URL https ://doi. org/10.1080/17442509608834055.
Stroock, D. W. Diffusion semigroups corresponding to uniformly elliptic divergence form operators. In Séminaire de Probabilités, XXII, volume 1321 of Lecture Notes in Math., pp. 316-347. Springer, Berlin (1988). doi:10.1007/BFb0084145. URL https://doi.org/10.1007/BFb0084145.
Trutnau, G. Multidimensional skew reflected diffusions. In Stochastic analysis: classical and quantum, pp. 228-244. World Sci. Publ., Hackensack, NJ (2005).


[^0]:    2010 Mathematics Subject Classification. 60G0, 60G12, 35K67,
    Key words and phrases. Probabilistic interpretation of divergence form operators, Dirichlet forms, skew diffusions, transmission problems.

    Miguel Martinez has been supported by the Labex Bézout.

