# BSDEs, càdlàg martingale problems and mean-variance hedging under basis risk 

Ismail Laachir, Francesco Russo

## To cite this version:

Ismail Laachir, Francesco Russo. BSDEs, càdlàg martingale problems and mean-variance hedging under basis risk. 2014. hal-01086227v1

## HAL Id: hal-01086227 <br> https://inria.hal.science/hal-01086227v1

Preprint submitted on 23 Nov 2014 (v1), last revised 20 Mar 2016 (v2)

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.

# BSDEs, càdlàg martingale problems and mean-variance hedging under basis risk 

Ismail LAACHIR ${ }^{(1,2)}$ and Francesco RUSSO $^{(1)}$

November 18th 2014


#### Abstract

.

The aim of this paper is to introduce a new formalism for the deterministic analysis associated with backward stochastic differential equations driven by general càdlàg martingales. When the martingale is a standard Brownian motion, the natural deterministic analysis is provided by the solution of a semilinear PDE of parabolic type. A significant application concerns the hedging problem under basis risk of a contingent claim $g\left(X_{T}, S_{T}\right)$, where $S$ (resp. $X$ ) is an underlying price of a traded (resp. non-traded but observable) asset, via the celebrated Föllmer-Schweizer decomposition. We revisit the case when the couple of price processes $(X, S)$ is a diffusion and we provide explicit expressions when $(X, S)$ is an exponential of additive processes. (1) ENSTA ParisTech, Unité de Mathématiques Appliquées, F-91120 Palaiseau, France. (2) Université de Bretagne Sud, Lab-STICC, F-56321 Lorient Cedex, France.

Key words and phrases: Backward stochastic differential equations; càdlàg martingales; basis risk; Föllmer-Schweizer decomposition; quadratic hedging; martingale problem.


MSC Classification 2010: 91G80; 60G51; 60J25; 60J75; 60H10; 60H30; 91G80

## 1 Introduction

The motivation of this work comes from the hedging problem in the presence of basis risk. When a derivative product is based on a non traded or illiquid underlying, the specification of a hedging strategy becomes problematic. In practice one frequent methodology consists in constituting a portfolio based on a (traded and liquid) additional asset which is correlated with the original one. The use of a non perfectly correlated asset induces a residual risk, often called basis risk, that makes the market incomplete. A common example is the hedging of
a basket (or index based) option, only using a subset of the assets composing the contract. Commodity markets also present many situations where basis risk plays an essential role, since many goods (as kerosene) do not have liquid future markets. For instance, kerosene consumers as airline companies, who want to hedge their exposure to the fuel use alternative future contracts, as crude oil or heating oil. The latter two commodities are strongly correlated to kerosene and their corresponding future market is liquid. Weather derivatives constitute an example of contract written on a non-traded underlying, since they are based on heating temperature; natural gas or electricity are in general used to hedge these contracts.

In this work, we consider a maturity $T>0$, a pair of processes $(X, S)$ and a contingent claim of the type $h:=g\left(X_{T}\right)$ or even $h:=g\left(X_{T}, S_{T}\right) . X$ is a non traded or illiquid, but observable asset and $S$ is a traded asset, correlated to $X$. In order to hedge this derivative, in general the practitioners use the proxy asset $S$ as a hedging instrument, but since the two assets are not perfectly correlated, a basis risk exists. Because of the incompleteness of this market, one should define a risk aversion criterion. One possibility is to use the utility function based approach to define the hedging strategy, see for example [9], [18], [25], [26]. Another approach is based on the quadratic hedging error criterion: it follows the idea of the seminal work of [13] that introduces the theoretical bases of the quadratic hedging in incomplete markets. In particular, they show the close relation of the quadratic hedging problem with a special semimartingale decomposition, known as the Föllmer-Schweizer (F-S) decomposition. The reader can consult [33, 34] for basic information on F-S decomposition, which provides the so called local risk minimizing hedging strategy and it is a significant tool for solving the mean variance hedging problem in an incomplete market.
[20] applied this general framework to the quadratic hedging under basis risk in a simple log-normal model. They consider for instance the two-dimensional Black-Scholes model for the non traded (but observable) $X$ and the hedging asset $S$, described by

$$
\begin{aligned}
d X_{t} & =\mu_{X} X_{t} d t+\sigma_{X} X_{t} d W_{t}^{X} \\
d S_{t} & =\mu_{S} S_{t} d t+\sigma_{S} S_{t} d W_{t}^{S}
\end{aligned}
$$

where $\left(W^{X}, W^{S}\right)$ is a standard correlated two-dimensional Brownian motion. They derive the F-S decomposition of a European payoff $h=g\left(X_{T}\right)$, i.e.

$$
\begin{equation*}
g\left(X_{T}\right)=h_{0}+\int_{0}^{t} Z_{s}^{h} d S_{s}+L_{T}^{h} \tag{1.1}
\end{equation*}
$$

where $L^{h}$ is a martingale which is strongly orthogonal to the martingale part of the hedging asset process $S$. Using the Feynman-Kac theorem, they relate the decomposition components $h_{0}$ and $Z^{h}$ to a PDE terminal-value problem. This yields, as byproduct, the price and hedging portfolio of the European option $h$. These quantities can be expressed in closed formulae in the case of call-put options. Extensions of those results to the case of stochastic correlation between the two assets $X$ and $S$, have been performed by [1].

Coming back to the general case, the $\mathrm{F}-\mathrm{S}$ decomposition of $h$ with respect to the $\mathcal{F}_{t^{-}}$ semimartingale $S$ can be seen as a special case of the well-known backward stochastic differential equations (BSDEs). We look for a triplet of processes $(Y, Z, O)$ being solution of an equation of the form

$$
\begin{equation*}
Y_{t}=h+\int_{t}^{T} \hat{f}\left(\omega, s, Y_{s-}, Z_{s}\right) d V_{s}^{S}-\int_{t}^{T} Z_{s} d M_{s}^{S}-\left(O_{T}-O_{t}\right) \tag{1.2}
\end{equation*}
$$

where $M^{S}$ (resp. $V^{S}$ ) is the local martingale (resp. the bounded variation process) appearing in the semimartingale decomposition of $S, O$ is a strongly orthogonal martingale to $M^{S}$, and $\hat{f}(\omega, s, y, z)=-z$.

BSDEs were first studied in the Brownian framework by [27] with an early paper of [3]. [27] showed existence and uniqueness of the solutions when the coefficient $\hat{f}$ is globally Lipschitz with respect to $(y, z)$ and $h$ being square integrable. It was followed by a long series of contributions, see for example [10] for a survey on Brownian based BSDEs and applications to finance. For example, the Lipschitz condition was essential in $z$ and only a monotonicity condition is required for $y$. Many other generalizations were considered. We also drive the attention on the recent monograph [28].

When the driving martingale in the BSDE is a Brownian motion, $h=g\left(S_{T}\right)$, and $S$ is a Markov diffusion, a solution of a BSDE constitutes a probabilistic representation of a semilinear parabolic PDE. In particular if $u$ is a solution of the mentioned PDE, then, roughly speaking setting $Y_{t}=u\left(t, S_{t}\right), Z=\partial_{s} u\left(t, S_{t}\right), O \equiv 0$, the triplet $(Y, Z, O)$ is a solution to (1.2). That PDE is a deterministic problem naturally related to the BSDE. When $M^{S}$ is a general càdlàg martingale, the link between a $\operatorname{BSDE}(1.2)$ and a deterministic problem is less obvious.

As far as backward stochastic differential equations driven by a martingale, the first paper seems to be [5]. Later, several authors have contributed to that subject, for instance [4] and [11]. More recently [6, Theorem 3.1] give sufficient conditions for existence and uniqueness for BSDEs of the form (1.2).

In this paper we consider a forward-backward SDE, issued from (1.2), where the forward process solves a sort of martingale problem (in the strong probability sense, i.e. where the underlying filtration is fixed) instead of the usual stochastic differential equation appearing in the Brownian case. More particularly we suppose the existence of an operator $\mathfrak{a}: \mathfrak{D}(\mathfrak{a}) \subset$ $C\left([0, T] \times \mathbb{R}^{2}\right) \rightarrow \mathcal{L}$, where $\mathcal{L}$ is a suitable space of functions $[0, T] \times \mathbb{R}^{2} \rightarrow \mathbb{C}^{2}($ see (2.3)), such that $(X, S)$ verifies the following:

$$
\forall g \in \mathfrak{D}(\mathfrak{a}), \quad\left(g\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \mathfrak{a}\left(u, X_{u-}, S_{u-}\right) d A_{u}\right)_{0 \leq t \leq T} \quad \text { is an } \mathcal{F}_{t} \text {-local martingale, }
$$

and $A$ is some fixed predictable bounded variation process. With $\mathfrak{a}$ we associate the operator $\widetilde{\mathfrak{a}}$ defined by

$$
\widetilde{\mathfrak{a}}(y):=\mathfrak{a}(\widetilde{y})-y \mathfrak{a}(i d)-i d \mathfrak{a}(y)
$$

where $i d(t, x, s)=s, \widetilde{y}=y \times i d$.
In the forward-backward BSDE we are interested in, the driver $\hat{f}$ verifies
$\mathfrak{a}(i d)\left(t, X_{t-}(\omega), S_{t-}(\omega)\right) \hat{f}(\omega, t, y, z)=f\left(t, X_{t-}(\omega), S_{t-}(\omega), y, z\right),(t, y, z) \in[0, T] \times \mathbb{C}^{2}, \omega \in \Omega$,
for some $f:[0, T] \times \mathbb{R}^{2} \times \mathbb{C}^{2} \rightarrow \mathbb{C}$. The main idea is to settle a deterministic problem which is naturally associated with the forward-backward SDE (1.2).

The deterministic problem consists in looking for a pair of functions $(y, z)$ which solves

$$
\begin{align*}
& \mathfrak{a}(y)(t, x, s)=-f(t, x, s, y(t, x, s), z(t, x, s)) \\
& \widetilde{\mathfrak{a}}(y)(t, x, s)=z(t, x, s) \widetilde{\mathfrak{a}}(i d)(t, x, s) \tag{1.4}
\end{align*}
$$

for all $t \in[0, T]$ and $(x, s) \in \mathbb{R}^{2}$, with the terminal condition $y(T, .,)=.g(.,$.$) .$

Any solution to the deterministic problem (1.4) will provide a solution $(Y, Z, O)$ to the corresponding BSDE, setting

$$
Y_{t}=y\left(t, X_{t}, S_{t}\right), Z_{t}=z\left(t, X_{t-}, S_{t-}\right)
$$

For illustration, let us consider the elementary case when $S$ is a diffusion process fulfilling $d S_{t}=\sigma_{S}\left(t, S_{t}\right) d W_{t}+b_{S}\left(t, S_{t}\right) d t$, and $X \equiv 0$. Then $A_{t} \equiv t,\left\langle M^{S}\right\rangle=\int_{0}^{2}\left(\sigma^{S}\right)^{2}\left(r, S_{r}\right) d r$, $V^{S}=\int_{0}^{r} b\left(r, S_{r}\right) d r=\int_{0}^{r} \mathfrak{a}(i d)\left(r, S_{r}\right) d r ; \mathfrak{a}$ is the parabolic generator of $S, \mathfrak{D}(\mathfrak{a})=C^{1,2}([0, T] \times$ $\mathbb{R}^{2} \rightarrow \mathbb{C}$. In that case (1.4) becomes

$$
\begin{align*}
\partial_{t} y(t, x, s)+\left(b_{s} \partial_{s} f+\frac{1}{2} \sigma_{S}^{2} \partial_{s s} f\right)(t, x, s) & =-f(t, x, s, y(t, x, s), z(t, x, s))  \tag{1.5}\\
z & =\partial_{s} y
\end{align*}
$$

In that situation $\tilde{\mathfrak{a}}$ is closely related to the classical derivation operator. When $S$ models the price of a traded asset and $f(t, x, s, y, z)=-\left(\sigma_{S}^{2}\right)(t, s) z$, the resolution of (1.5) emerging from the BSDE (1.2) with (1.3), allows to solve the usual (complete market Black-Scholes type) hedging problem with underlying $S$. Consequently, in the general case, $\tilde{\mathfrak{a}}$ appears to be naturally associated with a sort of "generalized derivation map". A first link between the hedging problem in incomplete markets and generalized derivation operators was observed for instance in [16].

The aim of our paper is threefold.

1) To provide a general methodology for solving forward-backward SDEs driven by a càdlàg martingale, via the solution of a deterministic problem generalizing the classical partial differential problem appearing in the case of Brownian martingales.
2) To give applications to the hedging problem in the case of basis risk via the FöllmerSchweizer decomposition. In particular we revisit the case when $(X, S)$ is a diffusion process whose particular case of Black-Scholes was treated by [20], discussing some analysis related to a corresponding PDE.
3) To furnish a quasi-explicit solution when the pair of processes $(X, S)$ is an exponential of additive processes, which constitutes a generalization of the results of [17] and [19], established in the absence of basis-risk. This yields a characterization of the hedging strategy in terms of Fourier-Laplace transform and the moment generating function.

The paper is organized as follows. In Section 2, we state the strong inhomogeneous martingale problem, and we give several examples, as Markov flows and the exponential of additive processes. In Section 3, we state the general form of a BSDE driven by a martingale and we associate a deterministic problem with it. We show in particular that a solution for this deterministic problem yields a solution for the BSDE. In Section 4, we apply previous methodology to the F-S decomposition problem under basis risk. In the case of exponential of additives processes, we obtain a quasi-explicit decomposition of the mentioned F-S decomposition.

## 2 Strong inhomogeneous martingale problem

### 2.1 General considerations

In this paper $T$ will be a strictly positive number. We consider a complete probability space $(\Omega, \mathcal{F}, \mathbb{P})$ with a filtration $\left(\mathcal{F}_{t}\right)_{t \in[0, T]}$, fulfilling the usual conditions. By default, all the processes will be indexed by $[0, T]$. Let $(X, S)$ a couple of $\mathcal{F}_{t}$-adapted processes. We will often mention concepts as martingale, semimartingale, adapted, predictable without mentioning the underlying filtration $\left(\mathcal{F}_{t}\right)_{t \in[0, T]}$. Given a bounded variation function $\phi:[0, T] \rightarrow \mathbb{R}$, we will denote by $t \mapsto\|\phi\|_{t}$ the associated total variation function.

We introduce a notion of martingale type problem related to $(X, S)$, which is a generalization of a stochastic differential equation. We emphasize that the present notion looks similar to the classical notion of [35] but here the notion is probabilistically strong and relies on a fixed filtered probability space. In the context of Stroock and Varadhan, however, a solution is a probability measure and the underlying process is the canonical process on some canonical space. Here a filtered probability space is given at the beginning. A similar notion was introduced in [31] Definition 5.1. A priori, we will not suppose that our strong martingale problem is well-posed (existence and uniqueness).

Definition 2.1. Let $\mathcal{O}$ be an open set of $\mathbb{R}^{2}$. Let $\left(A_{t}\right)$ be an $\mathcal{F}_{t}$-adapted bounded variation continuous process, such that, a.s.

$$
\begin{equation*}
d A_{t} \ll d \rho_{t} \tag{2.1}
\end{equation*}
$$

for some bounded variation function $\rho$, and $\mathfrak{a}$ a map

$$
\begin{equation*}
\mathfrak{a}: \mathfrak{D}(\mathfrak{a}) \subset \mathcal{C}([0, T] \times \mathcal{O}, \mathbb{C}) \longrightarrow \mathcal{L} \tag{2.2}
\end{equation*}
$$

where

$$
\begin{align*}
\mathcal{L}=\{f: & {[0, T] \times \mathcal{O} \rightarrow \mathbb{C}, \text { such that for every compact } K \text { of } \mathcal{O} } \\
& \left.\|f\|_{K}(t):=\sup _{(x, y) \in K}|f(t, x, y)|<\infty \quad d \rho_{t} \text { a.e. }\right\} \tag{2.3}
\end{align*}
$$

We say that a couple of càdlàg processes $(X, S)$ is a solution of the strong martingale problem related to $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$, if for any $g \in \mathfrak{D}(\mathfrak{a}),\left(g\left(t, X_{t}, S_{t}\right)\right)_{t}$ is a semimartingale with continuous bounded variation component such that

$$
\begin{equation*}
\int_{0}^{t}\left|\mathfrak{a}(g)\left(u, X_{u-}, S_{u-}\right)\right| d\|A\|_{u}<\infty \text { a.s. } \tag{2.4}
\end{equation*}
$$

and

$$
\begin{equation*}
t \longmapsto g\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(g)\left(u, X_{u-}, S_{u-}\right) d A_{u} \tag{2.5}
\end{equation*}
$$

is an $\mathcal{F}_{t^{-}}$local martingale.
We start introducing some significant notations.

## Notation 2.2.

1) $i d:(t, x, s) \longmapsto s$.
2) For any $y \in \mathcal{C}([0, T] \times \mathcal{O})$, we denote by $\widetilde{y}$ the function $\widetilde{y}:=y \times i d$, i.e.

$$
\begin{equation*}
(y \times i d)(t, x, s)=s y(t, x, s) \tag{2.6}
\end{equation*}
$$

3) Suppose that $i d \in \mathfrak{D}(\mathfrak{a})$. For $y \in \mathfrak{D}(\mathfrak{a})$ such that $\widetilde{y} \in \mathfrak{D}(\mathfrak{a})$, we set

$$
\begin{equation*}
\widetilde{\mathfrak{a}}(y):=\mathfrak{a}(\widetilde{y})-y \mathfrak{a}(i d)-i d \mathfrak{a}(y) \tag{2.7}
\end{equation*}
$$

As we have mentioned in the introduction, the map $\tilde{\mathfrak{a}}$ will play the role of a generalized derivative. We state first a preliminary lemma.

Lemma 2.3. Let $y \in \mathfrak{D}(\mathfrak{a})$. Suppose that $y, i d, y \times i d \in \mathfrak{D}(\mathfrak{a})$. We set $Y_{t}=y\left(t, S_{t}, X_{t}\right)$ and $M^{Y}$ be its martingale component given in (2.5). Then

$$
\left\langle M^{Y}, M^{S}\right\rangle_{t}=\int_{0}^{t} \widetilde{\mathfrak{a}}(y)\left(u, X_{u-}, S_{u-}\right) d A_{u}
$$

Proof. In order to compute the angle bracket $\left\langle M^{Y}, M^{S}\right\rangle$, we start by expressing the square bracket of $M^{Y}$ and $M^{S}$. First, note that, since $y, i d \in \mathfrak{D}(\mathfrak{a})$ and $A$ is a continuous process, then the bounded variation parts of the semimartingales $\left(S_{t}\right)_{t}$ and $\left(y\left(t, S_{t}, X_{t}\right)\right)_{t}$ are continuous. We have

$$
\begin{aligned}
{\left[M^{Y}, M^{S}\right]_{t} } & =[Y, S]_{t} \\
& =(S Y)_{t}-\int_{0}^{t} Y_{s-} d S_{s}-\int_{0}^{t} S_{s-} d Y_{s}
\end{aligned}
$$

where the first equality is justified by the fact that the square bracket of any process with a continuous bounded variation process vanishes. Using moreover the fact that $y \times i d \in \mathfrak{D}(\mathfrak{a})$, the process

$$
\begin{aligned}
{\left[M^{Y}, M^{S}\right]-\int_{0} \mathfrak{a}(y \times i d)\left(r, X_{r-}, S_{r-}\right) d A_{r}+\int_{0} y\left(r, X_{r-},\right.} & \left.S_{r-}\right) \mathfrak{a}(i d)\left(r, X_{r-}, S_{r-}\right) d A_{r} \\
& +\int_{0} S_{r-} \mathfrak{a}(y)\left(r, X_{r-}, S_{r-}\right) d A_{r}
\end{aligned}
$$

is an $\mathcal{F}_{t}$-local martingale.
Consequently, $\left[M^{Y}, M^{S}\right]$ is a special $\mathcal{F}_{t}$-semimartingale. Since $\left\langle M^{Y}, M^{S}\right\rangle-\left[M^{Y}, M^{S}\right]$ is a local martingale, the result follows by uniqueness of the decomposition of a special semimartingale.

In the sequel, we will make the following assumption.
Assumption 2.4. $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$ verifies the following axioms.
i) $i d \in \mathfrak{D}(\mathfrak{a})$.
ii) $(t, x, s) \mapsto s^{2} \in \mathfrak{D}(\mathfrak{a})$.

Corollary 2.5. Let $(X, S)$ be a solution of the strong martingale problem introduced in Definition 2.1 then, under Assumption 2.4, $S$ is a special semimartingale with decomposition $M^{S}+V^{S}$ given below.
i) $V_{t}^{S}=\int_{0}^{t} \mathfrak{a}(i d)\left(u, X_{u-}, S_{u-}\right) d A_{u}$.
ii) $\left\langle M^{S}\right\rangle_{t}=\int_{0}^{t} \widetilde{\mathfrak{a}}(i d)\left(u, X_{u-}, S_{u-}\right) d A_{u}$.

Proof. i) is obvious since $i d \in \mathfrak{D}(\mathfrak{a})$ and ii) is a consequence of Lemma 2.3 and the fact that $(t, x, s) \mapsto s^{2}$ belongs to $\mathfrak{D}(\mathfrak{a})$.

In many situations, the operator $\mathfrak{a}$ is related to the classical infinitesimal generator, when it exists. We will make this relation explicit in the below example of Markov processes.

### 2.2 The case of Markov semigroup

In this section $\mathcal{O}$ will be for simplicity $\mathbb{R}^{2}$. In this example, for illustration, we only consider a single process $S$ instead of a couple $(X, S)$. For this reason, it is more comfortable to reexpress Definition 2.1 into the following simplified version. In the present case we will always have $A_{t} \equiv t$.

Definition 2.6. We say that $S$ is a solution of the strong martingale problem related to $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$, if there is a map

$$
\begin{equation*}
\mathfrak{a}: \mathfrak{D}(\mathfrak{a}) \subset \mathcal{C}([0, T] \times \mathbb{R}) \longrightarrow \mathcal{L} \tag{2.8}
\end{equation*}
$$

where

$$
\begin{align*}
\mathcal{L}=\{f: & {[0, T] \times \mathbb{R} \rightarrow \mathbb{C}, \text { such that for every compact } K \text { of } \mathbb{R} } \\
& \left.\|f\|_{K}(t):=\sup _{x \in K}|f(t, x)|<\infty \quad \text { dt a.e. }\right\}, \tag{2.9}
\end{align*}
$$

such that for any $g \in \mathfrak{D}(\mathfrak{a}),\left(g\left(t, S_{t}\right)\right)_{t}$ is a (special) semimartingale with continuous bounded variation component verifying

$$
\begin{equation*}
\int_{0}^{t}\left|\mathfrak{a}(g)\left(u, S_{u-}\right)\right| d u<\infty \text { a.s. } \tag{2.10}
\end{equation*}
$$

and

$$
\begin{equation*}
t \longmapsto g\left(t, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(g)\left(u, S_{u-}\right) d u \tag{2.11}
\end{equation*}
$$

is a $\mathcal{F}_{t}^{S}$ - local martingale, where $\mathcal{F}_{t}^{S}$ is the canonical filtration associated with $S$.
Let $\left(X_{t}^{s, x}\right)_{t \geq s, x \in \mathbb{R}}$ be a time-homogeneous Markovian flow. In particular, if $S=X^{\ell, x}$ and $f$ is a bounded Borel function, then

$$
\begin{equation*}
\mathbb{E}\left[f\left(S_{t}\right) \mid \mathcal{F}_{s}^{S}\right]=\Psi\left(t-s, S_{s}\right), \tag{2.12}
\end{equation*}
$$

where $\ell \leq s \leq t \leq T$ and

$$
\begin{equation*}
\Psi(r, y)=\mathbb{E}\left[f\left(X_{r}^{0, y}\right)\right]=\mathbb{E}\left[f\left(X_{s+r}^{s, y}\right)\right], \tag{2.13}
\end{equation*}
$$

for any $r, s \geq 0$ and $\mathcal{F}^{S}$ is the canonical filtration for $S$. We suppose moreover that $X_{t}^{\ell, x}$ is square integrable for any $0 \leq \ell \leq t \leq T$ and $x \in \mathbb{R}$. We denote by $E$ the linear space of functions such that

$$
\begin{equation*}
E=\left\{f \in C \text { such that } \tilde{f}:=x \mapsto \frac{f(x)}{1+x^{2}} \text { is uniformly continuous and bounded }\right\}, \tag{2.14}
\end{equation*}
$$

equipped with the norm

$$
\|f\|_{E}:=\sup _{x} \frac{|f(x)|}{1+x^{2}}<\infty .
$$

The set $E$ can easily be shown to be a Banach space equipped with the norm $\|\cdot\|_{E}$. Indeed $E$ is a suitable space for Markov processes which are square integrable. In particular, (2.12) and (2.13) remain valid if $f \in E$. From now on we consider the family of linear operators ( $P_{t}, t \geq 0$ ) defined on the space $E$ by

$$
\begin{equation*}
P_{t} f(x)=\mathbb{E}\left[f\left(X_{t}^{0, x}\right)\right], \text { for } t \in[0, T], x \in \mathbb{R}, \quad \forall f \in E . \tag{2.15}
\end{equation*}
$$

We formulate now a fundamental assumption.

## Assumption 2.7.

i) $P_{t} E \subset E$ for all $t \in[0, T]$.
ii) The linear operator $P_{t}$ is bounded, for all $t \in[0, T]$.
iii) $\left(P_{t}\right)$ is strongly continuous, i.e. $\lim _{t \rightarrow 0} P_{t} f=f$ in the $E$ topology.

Using the Markov flow property (2.12), it is easy to see that the family of continuous operators $\left(P_{t}\right)$ defined above has the semigroup property. In particular, under Assumption 2.7 , the family $\left(P_{t}\right)$ is strongly continuous semigroup on $E$.

Assumption 2.7 is fulfilled in many common cases, as mentioned in Proposition 2.8 and Remarks 2.9 and 2.10.

The proposition below concerns the validity of items i) and ii).
Proposition 2.8. Let $t \in[0, T]$. Suppose that $x \mapsto X_{t}^{0, x}$ is differentiable in $L^{2}(\Omega)$ such that

$$
\begin{equation*}
\sup _{x \in \mathbb{R}} \mathbb{E}\left[\left|\partial_{x} X_{t}^{0, x}\right|^{2}\right]<\infty \tag{2.16}
\end{equation*}
$$

Then $P_{t} f \in E$ for all $f \in E$ and $P_{t}$ is a bounded operator.
The proof of this proposition is reported in Appendix A.
Remark 2.9. Condition (2.16) of Proposition 2.8 is fulfilled in the following two cases.

1) If $\left(L_{t}\right)$ is a Lévy process, the Markov flow $X^{0, x}=x+L$ verifies $\partial_{x} X^{0, x}=1$.
2) If $\left(X_{t}^{0, x}\right)$ is a diffusion process verifying

$$
X_{t}^{0, x}=x+\int_{0}^{t} b\left(X_{s}^{0, x}\right) d s+\int_{0}^{t} \sigma\left(X_{s}^{0, x}\right) d W_{s}
$$

where $b$ and $\sigma$ are $C_{b}^{1}$ functions.
Remark 2.10. Item iii) of Assumption 2.7 is verified in the case of square integrable Lévy processes, c.f. Proposition B. 1 in Appendix B.

For the rest of this subsection we work under Assumption 2.7.
Item iii) of Assumption 2.7 permits to introduce the definition of the generator of $\left(P_{t}\right)$ as follows.

Definition 2.11. The generator $L$ of $\left(P_{t}\right)$ in $E$ is defined on the domain $D(L)$ which is the subspace of $E$ defined by

$$
\begin{equation*}
D(L)=\left\{f \in E \text { such that } \lim _{t \rightarrow 0} \frac{P_{t} f-f}{t} \text { exists in } E\right\} \tag{2.17}
\end{equation*}
$$

We denote by Lf the limit above. We refer to [21, Chapter 4], for more details.
Remark 2.12. If $f \in E$ such that there is $g \in E$ such that

$$
\left(P_{t} f\right)(x)-f(x)-\int_{0}^{t} P_{s} g(x) d s=0, \forall t \geq 0, x \in E
$$

then $f \in D(L)$ and $g=L f$.
Previous integral is always defined as E-valued Bochner integral. Indeed, since $\left(P_{t}\right)$ is strongly continuous, then by [21, Lemma 4.1.7], we have

$$
\begin{equation*}
\left\|P_{t}\right\| \leq M_{w} e^{w t} \tag{2.18}
\end{equation*}
$$

for some real $w$ and related constant $M_{w} \cdot\|\cdot\|$ denotes here the operator norm.
A useful result which allows to deal with time-dependent functions is given below.
Lemma 2.13. Let $f:[0, T] \rightarrow D(L) \subset E$. We suppose the following.
i) $f$ is continuous as a $D(L)$-valued function, where $D(L)$ is equipped with the graph norm.
ii) $f:[0, T] \rightarrow E$ is of class $C^{1}$.

Then, the below E-valued equality holds:

$$
\begin{equation*}
P_{t} f(t, .)=f(0, .)+\int_{0}^{t} P_{s}(L f(s, .)) d s+\int_{0}^{t} P_{s}\left(\frac{\partial f}{\partial s}(s, .)\right) d s, \forall t \in[0, T] \tag{2.19}
\end{equation*}
$$

Remark 2.14. We observe that the two integrals above can be considered as E-valued Bochner integrals because, by hypothesis, $s \mapsto L f(s, \cdot)$ and $s \mapsto \frac{\partial f}{\partial s}(s,$.$) are continuous with$ values in E, and so we can apply again (2.18) in Remark 2.12.

Proof. It will be enough to show that

$$
\begin{equation*}
\frac{d}{d t} P_{t} f(t, .)=P_{t}(L f(t, .))+P_{t}\left(\frac{\partial f}{\partial t}(t, .)\right), \forall t \in[0, T] . \tag{2.20}
\end{equation*}
$$

In fact, even if Banach space valued, a differentiable function at each point is also absolutely continuous.

Since the right-hand side of (2.20) is continuous it is enough to show that the rightderivative of $t \mapsto P_{t} f(t, \cdot)$ coincides with the right-hand side of (2.20). Let $h>0$. We evaluate

$$
P_{t+h} f(t+h, .)-P_{t} f(t, .)=I_{1}(t, h)+I_{2}(t, h)
$$

where

$$
I_{1}(t, h)=P_{t+h} f(t+h, .)-P_{t} f(t+h, .), I_{2}(t, h)=P_{t} f(t+h, .)-P_{t} f(t, .)
$$

Now by [21, Lemma 4.1.14], we get

$$
I_{1}(t, h):=P_{t+h} f(t+h, .)-P_{t} f(t+h, .)=\int_{t}^{t+h} P_{s} L f(t+h, .) d s
$$

We divide by $h$ and we get

$$
\begin{aligned}
\left\|\frac{1}{h} \int_{t}^{t+h}\left(P_{s} L f(t+h, .)-P_{s} L f(t, .)\right) d s\right\|_{E} & \leq \frac{1}{h} \int_{t}^{t+h} d s\left\|P_{s}\{L f(t+h, .)-L f(t, .)\}\right\|_{E} \\
& \leq\|L f(t+h, .)-L f(t, .)\|_{E} \frac{1}{h} \int_{t}^{t+h}\left\|P_{s}\right\| d s \\
& \leq\|f(t+h, .)-f(t, .)\|_{D(L)} \frac{1}{h} \int_{t}^{t+h}\left\|P_{s}\right\| d s
\end{aligned}
$$

where $\|\cdot\|_{D(L)}$ is the graph norm of $L$. This converges to zero (note that $\left\|P_{s}\right\|$ is bounded by (2.18)), and we get that

$$
\frac{1}{h} I_{1}(t, h) \xrightarrow{h \rightarrow 0} P_{t}(L f(t, .)) .
$$

We estimate now $I_{2}(t, h)$.

$$
\left\|\frac{P_{t} f(t+h, .)-P_{t} f(t, .)}{h}-P_{t}\left(\frac{\partial f}{\partial t}(t, .)\right)\right\|_{E} \leq\left\|P_{t}\right\|\left\|\frac{f(t+h, .)-f(t, .)}{h}-\frac{\partial f}{\partial t}(t, .)\right\|_{E} .
$$

This goes to zero as $h$ goes to zero, by Assumption ii).
This concludes the proof of Lemma 2.13.

We can now discuss the fact that a process $S=X^{0, x}$, where $X_{t}^{s, x}$ is a Markovian flow (as precised at the beginning of Section 2.2) is a solution to our (time inhomogeneous) strong martingale problem (2.6).

Theorem 2.15. We denote
$\mathfrak{D}(\mathfrak{a})=\{g:[0, T] \rightarrow D(L)$ such that assumptions i) and ii) of Lemma 2.13 are fulfilled $\}$
and for $g \in \mathfrak{D}(\mathfrak{a})$

$$
\mathfrak{a}(g)(t, x)=\frac{\partial g}{\partial t}(t, x)+L g(t, \cdot)(x), \forall t \in[0, T], x \in \mathbb{R} .
$$

Then $S$ is a solution of the strong martingale problem introduced in Definition 2.6.
Remark 2.16. Let $g \in \mathfrak{D}(\mathfrak{a})$. Since for each $t \in[0, T]$, by assumptions i) and ii) of Lemma 2.13, $\mathfrak{a}(g)(t, \cdot) \in E$, then, obviously $\mathfrak{a}(g) \in \mathcal{L}$. Moreover, the same assumptions imply that $t \mapsto \frac{\partial g}{\partial t}(t, \cdot)$ and $t \mapsto L g(t, \cdot)$ are continuous on $[0, T]$ and hence are bounded, i.e.

$$
\sup _{t \in[0, T]}\left\|\frac{\partial g}{\partial t}(t, \cdot)\right\|_{E}<\infty, \quad \sup _{t \in[0, T]}\|L g(t, \cdot)\|_{E}<\infty .
$$

This yields in particular that Condition (2.10) is fulfilled.
Proof of Theorem 2.15.
It remains to show the martingale property (2.11). We fix $0 \leq s<t \leq T$ and a bounded random variable $\mathcal{F}_{s}^{S}$-measurable $G$. It will be sufficient to show that

$$
\begin{equation*}
\mathbb{E}[A(s, t)]=0, \tag{2.21}
\end{equation*}
$$

where

$$
\begin{equation*}
A(s, t)=G\left(g\left(t, S_{t}\right)-g\left(s, S_{s}\right)-\int_{s}^{t} \partial_{r} g\left(r, S_{r}\right) d r-\int_{s}^{t} L g(r, .)\left(S_{r}\right) d r\right) . \tag{2.22}
\end{equation*}
$$

By taking the conditional expectation of $A(s, t)$ with respect to $\mathcal{F}_{s}^{S}$, using (2.12) and Fubini's theorem, we get

$$
\mathbb{E}\left[A(s, t) \mid \mathcal{F}_{s}^{S}\right]=G \phi\left(S_{s}\right),
$$

where
$\phi(x)=\left(P_{t-s} g(t,).(x)-g(s, x)-\int_{s}^{t}\left(P_{r-s} \partial_{r} g(r,).\right)(x) d r-\int_{s}^{t}\left(P_{r-s} L g(r,).\right)(x) d r\right), \forall x \in \mathbb{R}$.
We define $f:[0, T-s] \times \mathbb{R} \rightarrow \mathbb{R}$ by $f(\tau, \cdot)=g(\tau+s, \cdot)$. $f$ fulfills the assumptions of Lemma 2.13 with $T$ being replaced by $T-s$. By the change of variable $u=r-s$, setting $\tau=t-s$, the equality above becomes

$$
\phi(x)=\left(P_{\tau} f(\tau, .)(x)-f(0, x)-\int_{0}^{\tau}\left(P_{u} \partial_{r} f(u, .)\right)(x) d u-\int_{0}^{\tau}\left(P_{u} L f(u, .)\right)(x) d r\right) .
$$

Now by Lemma 2.13 we get that $\phi(x)=0, \forall x \in \mathbb{R}$. Consequently $\mathbb{E}\left[A(s, t) \mid \mathcal{F}_{s}^{S}\right]=0$ and (2.21) is fulfilled.

Remark 2.17. We introduce the following subspace $E_{0}^{2}$ of $C^{2}$.

$$
\begin{equation*}
E_{0}^{2}=\left\{f \in C^{2} \text { such that } f^{\prime \prime} \text { vanishes at infinity }\right\} . \tag{2.23}
\end{equation*}
$$

Note that only the second derivative is supposed to vanish at infinity.
$E_{0}^{2}$ is included in $E$. Indeed, if $f \in E_{0}^{2}$, then the Taylor expansion $f(x)=f(0)+x f^{\prime}(0)+$ $x^{2} \int_{0}^{1}(1-\alpha) f^{\prime \prime}(x \alpha) d \alpha$ implies that $\tilde{f}$ is bounded. On the other hand, by straightforward calculus we see that the first derivative $\frac{d \tilde{f}}{d x}$ is bounded. This implies that $\tilde{f}$ is uniformly continuous.
In several examples it is easy to identify $E_{0}^{2}$ as a significant subspace of $D(L)$, see for instance the example of Lévy processes which is described below.

### 2.2.1 A significant particular case: Lévy processes

As anticipated above, an insightful example for Markov flows is the case of Lévy processes. We specify below the corresponding infinitesimal generator.

Let $\left(X_{t}\right)$ be a square integrable Lévy process with characteristic triplet $(A, \nu, \gamma)$, such that $X_{0}=0$. We refer to, e.g., $[8$, Chapter 3] for more details.

We suppose that $\left(X_{t}\right)$ is of pure jump, i.e. $A=0$ and $\gamma=0$. Since $X$ is square integrable, then (c.f. [8, Proposition 3.13])

$$
\begin{equation*}
\int_{\mathbb{R}}|x|^{2} \nu(d x)<\infty \tag{2.24}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{1}:=\frac{\mathbb{E}\left[X_{t}\right]}{t}=\int_{|x|>1} x \nu(d x)<\infty, \quad c_{2}:=\frac{\operatorname{Var}\left[X_{t}\right]}{t}=\int_{\mathbb{R}}|x|^{2} \nu(d x)<\infty . \tag{2.25}
\end{equation*}
$$

Clearly the corresponding Markov flow is given by $X_{t}^{0, x}=x+X_{t}, t \geq 0, x \in \mathbb{R}$.
The classical theory of semigroup for Lévy processes is for instance developed in Section 6.31 of [32]. There one defines the semigroup $P$ on the set $C_{0}$ of continuous functions vanishing at infinity, equipped with the sup-norm $\|u\|_{\infty}=\sup _{x}|u(x)|$. By [32, Theorem 31.5], the semigroup $P$ is strongly continuous on $C_{0}$, with norm $\|P\|=1$, and its generator $L_{0}$ is given by

$$
\begin{equation*}
L_{0} f(x)=\int\left(f(x+y)-f(x)-y f^{\prime}(x) 1_{|y|<1}\right) \nu(d y) \tag{2.26}
\end{equation*}
$$

Moreover, the set $C_{0}^{2}$ of functions $f \in C^{2}$ such that $f, f^{\prime}$ and $f^{\prime \prime}$ vanish at infinity, is included in $D\left(L_{0}\right)$. We remark that the domain $D(L)$ includes the classical domain $D\left(L_{0}\right)$. In fact, we have

$$
\|g\|_{E} \leq\|g\|_{C_{0}}, \forall g \in C_{0}
$$

Consequently, if $f \in D\left(L_{0}\right) \subset C_{0}$, then for $t>0$

$$
\left\|\frac{P_{t} f-f}{t}-L_{0} f\right\|_{E} \leq\left\|\frac{P_{t} f-f}{t}-L_{0} f\right\|_{C_{0}}
$$

So $f \in D(L)$ and $L f=L_{0} f$.
Assumption 2.7 is verified because of Proposition 2.8, item 1) of Remark 2.9 and Remark 2.10 .

The theorem below shows that the space $E_{0}^{2}$, defined in Remark 2.17, is a subset of $D(L)$.
Theorem 2.18. Let $L$ be the infinitesimal generator of the semigroup $\left(P_{t}\right)$. Then $E_{0}^{2} \subset D(L)$ and

$$
\begin{equation*}
L f(x)=\int\left(f(x+y)-f(x)-y f^{\prime}(x) 1_{|y|<1}\right) \nu(d y), f \in E_{0}^{2} \tag{2.27}
\end{equation*}
$$

A proof of this result, using arguments in [12], is developed in Appendix B.

In conclusion, $C^{2}$ functions whose second derivative vanishes at infinity belong to $D(L)$. For instance, id: $x \mapsto x \in D(L)$. On the other hand the function $x^{2}: x \mapsto x^{2}$ also belongs to $D(L)$.

In fact, for every $x \in \mathbb{R}, t \geq 0$ we have

$$
P_{t} f(x)-f(x)=\frac{\mathbb{E}\left[\left(x+X_{t}\right)^{2}\right]-x^{2}}{t}=\frac{2 x c_{1} t+c_{2} t+c_{1}^{2} t^{2}}{t}
$$

Obviously, this converges to the function $x \mapsto 2 x c_{1}+c_{2}$ according to the $E$-norm. Finally it follows that $L\left(x^{2}\right)=2 x c_{1}+c_{2}$.

Corollary 2.19. We have the following inclusion:

$$
E_{0}^{2} \oplus x^{2} \subset D(L)
$$

### 2.3 Diffusion processes

Here we will suppose again $\mathcal{O}=\mathbb{R} \times E$, where $E=\mathbb{R}$ or $] 0, \infty[$. A function $f:[0, T] \times \mathcal{O}$ will be said to be globally Lipschitz if it is Lipschitz with respect to the second and third variable uniformly with respect to the first.

We consider here the case of a diffusion process $(X, S)$ whose dynamics is described as follows:

$$
\begin{align*}
d X_{t} & =b_{X}\left(t, X_{t}, S_{t}\right) d t+\sum_{i=1}^{2} \sigma_{X, i}\left(t, X_{t}, S_{t}\right) d W_{t}^{i}  \tag{2.28}\\
d S_{t} & =b_{S}\left(t, X_{t}, S_{t}\right) d t+\sum_{i=1}^{d} \sigma_{S, i}\left(t, X_{t}, S_{t}\right) d W_{t}^{i}
\end{align*}
$$

where $W=\left(W^{1}, W^{2}\right)$ is a standard 2-dimensional Brownian motion with canonical filtration $\left(\mathcal{F}_{t}\right), b_{X}, b_{S}, \sigma_{X, i}$, and $\sigma_{S, i}$, for $i=1,2, b, \sigma:[0, T] \times \mathbb{R}^{2} \rightarrow \mathbb{R}$ are continuous functions which are globally Lipschitz.

We also suppose $\left(X_{0}, S_{0}\right)$ to have all moments and that $S$ takes value in $E$. For instance a geometric Brownian motion takes value in $E=] 0, \infty[$, if it starts from a positive point.

Remark 2.20. Let $p \geq 1$. It is well-known, that there is a constant $C(p)$, only depending on p, such that

$$
\mathbb{E}\left[\sup _{t \leq T}\left(\left|X_{t}\right|^{p}+\left|S_{t}\right|^{p}\right)\right] \leq C(p)\left(\left|X_{0}\right|^{p}+\left|S_{0}\right|^{p}\right)
$$

By Itô formula, for $f \in \mathcal{C}^{1,2}([0, T[\times \mathcal{O})$, we have

$$
\begin{aligned}
d f\left(t, X_{t}, S_{t}\right) & =\partial_{t} f\left(t, X_{t}, S_{t}\right) d t+\partial_{s} f\left(t, X_{t}, S_{t}\right) d S_{t}+\partial_{x} f\left(t, X_{t}, S_{t}\right) d X_{t} \\
& +\frac{1}{2}\left\{\partial_{s s} f\left(t, X_{t}, S_{t}\right) d\langle S\rangle_{t}+\partial_{x x} f\left(t, S_{t}, X_{t}\right) d\langle X\rangle_{t}+2 \partial_{s x} f\left(t, X_{t}, S_{t}\right) d\langle S, X\rangle_{t}\right\}
\end{aligned}
$$

We denote $\left|\sigma_{S}\right|^{2}=\sum_{i=1}^{2} \sigma_{S, i}^{2},\left|\sigma_{X}\right|^{2}=\sum_{i=1}^{2} \sigma_{X, i}^{2}$ and $\left\langle\sigma_{S}, \sigma_{X}\right\rangle=\sum_{i=1}^{2} \sigma_{S, i} \sigma_{X, i}$.
Hence, the operator $\mathfrak{a}$ can be defined as

$$
\begin{aligned}
\mathfrak{a}(f) & =\partial_{t} f+b_{S} \partial_{s} f+b_{X} \partial_{x} f \\
& +\frac{1}{2}\left\{\left|\sigma_{S}\right|^{2} \partial_{s s} f+\left|\sigma_{X}\right|^{2} \partial_{x x} f+2\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{s x} f\right\}
\end{aligned}
$$

associated with $A_{t} \equiv t$ and a domain $\mathfrak{D}(\mathfrak{a})=\mathcal{C}^{1,2}\left(\left[0, T[\times \mathcal{O}) \cap \mathcal{C}^{0}([0, T] \times \mathcal{O})\right.\right.$.
Note that Assumption 2.4 is verified since $i d$ and $i d \times i d$ belong to $\mathfrak{D}(\mathfrak{a})$. Moreover, a straightforward calculation gives

$$
\tilde{\mathfrak{a}}(f)=\left|\sigma_{S}\right|^{2} \partial_{s} f(t, x, s)+\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{x} f(t, x, s)
$$

In particular,

$$
\widetilde{\mathfrak{a}}(i d)=\left|\sigma_{S}\right|^{2}
$$

Remark 2.21. By Itô formula, for $0 \leq u \leq T$, we obviously have

$$
\begin{aligned}
f\left(u, X_{u}, S_{u}\right)-\int_{0}^{u} \mathfrak{a}(f)\left(r, X_{r}, S_{r}\right) d r & =\int_{0}^{u} \partial_{x} f\left(r, X_{r}, S_{r}\right)\left(\sigma_{X, 1}\left(r, X_{r}, S_{r}\right) d W_{r}^{1}+\sigma_{X, 2}\left(r, X_{r}, S_{r}\right) d W_{r}^{2}\right) \\
& +\int_{0}^{u} \partial_{s} f\left(r, X_{r}, S_{r}\right)\left(\sigma_{S, 1}\left(r, X_{r}, S_{r}\right) d W_{r}^{1}+\sigma_{S, 2}\left(r, X_{r}, S_{r}\right) d W_{r}^{2}\right)
\end{aligned}
$$

### 2.4 Variant of diffusion processes

Let $\left(W_{t}\right)$ be an $\mathcal{F}_{t}$-standard Brownian motion and $S$ be a solution of the SDE

$$
\begin{equation*}
d S_{t}=\sigma\left(t, S_{t}\right) d W_{t}+b_{1}\left(t, S_{t}\right) d a_{t}+b_{2}\left(t, S_{t}\right) d t \tag{2.29}
\end{equation*}
$$

where $b_{1}, b_{2}, \sigma:[0, T] \times \mathbb{R}^{2} \rightarrow \mathbb{R}$ are continuous functions which are globally Lipschitz, and $a:[0, T] \rightarrow \mathbb{R}$ is an increasing function such that $d a$ is singular with respect to Lebesgue measure. We set $A_{t}=a_{t}+t$.

The equation (2.29) can be written as

$$
d S_{t}=\sigma\left(t, S_{t}\right) d W_{t}+\left(b_{1}\left(t, S_{t}\right) \frac{d \rho_{t}}{d A_{t}}+b_{2}\left(t, S_{t}\right) \frac{d t}{d A_{t}}\right) d A_{t}
$$

A solution $S$ of (2.29) verifies the strong martingale problem with respect to ( $\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$, in the sense where $\mathfrak{D}(\mathfrak{a})=\mathcal{C}^{1,2}([0, T] \times \mathbb{R})$ and for $f \in \mathfrak{D}(\mathfrak{a})$,

$$
\mathfrak{a}(f)(t, s)=\left(\partial_{t} f(r, s) \frac{d r}{d A_{r}}+\partial_{s} f(r, s) \widetilde{b}(r, s)+\frac{1}{2} \partial_{s s} f(r, s) \widetilde{\sigma}^{2}(r, s)\right),
$$

where $\widetilde{b}(t, s)=b_{1}(t, s) \frac{d \rho_{t}}{d A_{t}}(t)+b_{2}(t, s) \frac{d t}{d A_{t}}(t)$ and $\widetilde{\sigma}^{2}(t, s)=\sigma^{2}(t, s) \frac{d t}{d A_{t}}(t)$.
Indeed, by Itô formula, the process

$$
t \mapsto f\left(t, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(f)\left(r, S_{r}\right) d A_{r}
$$

is a local martingale.

### 2.5 Exponential of additive processes

A càdlàg process $\left(Z^{1}, Z^{2}\right)$ is said to be an additive process if $\left(Z^{1}, Z^{2}\right)_{0}=0,\left(Z^{1}, Z^{2}\right)$ is continuous in probability and it has independent increments, i.e. $\left(Z_{t}^{1}-Z_{s}^{1}, Z_{t}^{2}-Z_{s}^{2}\right)$ is independent of $\mathcal{F}_{s}$ for $0 \leq s \leq t \leq T$ and $\left(\mathcal{F}_{s}\right)$ is the canonical filtration associated with $\left(Z^{1}, Z^{2}\right)$.

In this section we restrict ourselves to the case of exponential of additive processes which are semimartingales (shortly semimartingale additive processes) and we specify a corresponding martingale problem $(\mathfrak{a}, \mathfrak{D}(\mathfrak{a}), A)$ for this process. This will be based on Fourier-Laplace transform techniques. The couple of processes $(X, S)$ is defined by

$$
\begin{aligned}
X & =\exp \left(Z^{1}\right) \\
S & =\exp \left(Z^{2}\right)
\end{aligned}
$$

where $\left(Z^{1}, Z^{2}\right)$ is an (two-dimensional) semimartingale additive process.
We denote by $D$ the set

$$
D:=\left\{z=\left(z_{1}, z_{2}\right) \in \mathbb{C}^{2} \mid \quad \mathbb{E}\left[\left|X_{T}^{\operatorname{Re}\left(z_{1}\right)} S_{T}^{\operatorname{Re}\left(z_{2}\right)}\right|\right]<\infty\right\}
$$

We convene that $\mathbb{C}^{2}=\mathbb{R}^{2}+i \mathbb{R}^{2}$, associating the couple $\left(z_{1}, z_{2}\right)$ with $\left(\operatorname{Re} z_{1}, \operatorname{Re} z_{2}\right)+i\left(\operatorname{Im} z_{1}, \operatorname{Im} z_{2}\right)$. Clearly we have $D=\left(D \cap \mathbb{R}^{2}\right)+i \mathbb{R}^{2}$. We also introduce the notation

$$
D / 2:=\left\{z \in \mathbb{C}^{2} \mid \quad 2 z \in D\right\} \subset D
$$

Remark 2.22. By Cauchy-Schwarz inequality, $z, y \in D / 2$ implies that $z+y \in D$.
We denote by $\kappa: D \rightarrow \mathbb{C}$, the generating function of $\left(Z^{1}, Z^{2}\right)$, see for instance $[17$, Definition 2.1]. In particular $\kappa$ verifies

$$
\exp \left(\kappa_{t}\left(z_{1}, z_{2}\right)\right)=\mathbb{E}\left[\exp \left(z_{1} Z_{t}^{1}+z_{2} Z_{t}^{2}\right)\right]=\mathbb{E}\left[X_{t}^{z_{1}} S_{t}^{z_{2}}\right]
$$

We will adopt similar notations and assumptions as in [17], which treated the problem of variance optimal hedging for a one-dimensional exponential of additive process. We introduce a function $\rho$, defined, for each $t \in[0, T]$, as follows:

$$
\begin{align*}
\rho_{t}\left(z_{1}, z_{2}, y_{1}, y_{2}\right) & :=\kappa_{t}\left(z_{1}+y_{1}, z_{2}+y_{2}\right)-\kappa_{t}\left(z_{1}, z_{2}\right)-\kappa_{t}\left(y_{1}, y_{2}\right), \quad \text { for }\left(z_{1}, z_{2}\right),\left(y_{1}, y_{2}\right) \in D / 2 \\
\rho_{t}\left(y_{1}, y_{2}\right) & :=\rho_{t}\left(z_{1}, z_{2}, \overline{z_{1}}, \overline{z_{2}}\right), \quad \text { for }\left(z_{1}, z_{2}\right) \in D / 2 \\
\rho_{t}^{S} & :=\rho_{t}(0,1)=\kappa_{t}(0,2)-2 \kappa_{t}(0,1), \text { if }(0,1) \in D / 2 \tag{2.30}
\end{align*}
$$

We remark that for $\left(z_{1}, z_{2}\right) \in D / 2, t \mapsto \rho_{t}\left(z_{1}, z_{2}\right)$ is a real function. These functions appear naturally in the expression of the angle brackets of $\left(M^{X}, M^{S}\right)$ where $M^{X}$ (resp. $M^{S}$ ) is the martingale part of $X$ (resp. $S$ ).

From now on, in this section, the assumption below will be in force.

## Assumption 2.23.

1) $\rho^{S}$ is strictly increasing.
2) $(0,2) \in D$. This is equivalent to the existence of the second order moment of $S$.

Note that, by Cauchy-Schwarz, the second item implies that, $D / 2+(0,1) \subset D$.
We remind that previous assumption implies that $Z^{2}$ has no deterministic increments, see [17, Lemma 3.9].

Similarly as in [17, Propositions 3.4 and 3.15$]$, one can prove the following.

## Proposition 2.24.

1) For every $\left(z_{1}, z_{2}\right) \in D,\left(X_{t}^{z_{1}} S_{t}^{z_{2}} e^{-\kappa_{t}\left(z_{1}, z_{2}\right)}\right)$ is a martingale.
2) $t \mapsto \kappa_{t}\left(z_{1}, z_{2}\right)$ is a bounded variation continuous function, for every $\left(z_{1}, z_{2}\right) \in D$. In particular, $t \mapsto \rho_{t}\left(z_{1}, z_{2}\right)$ is also a bounded variation continuous function, for every $\left(z_{1}, z_{2}\right) \in D / 2$.
3) Let I be a compact real set included in $D$. Then

$$
\sup _{(x, y) \in I} \sup _{t \leq T} \mathbb{E}\left[X_{t}^{x} S_{t}^{y}\right]=\sup _{(x, y) \in I} \sup _{t \leq T} e^{\kappa_{t}(x, y)}<\infty
$$

4) $\forall\left(z_{1}, z_{2}\right) \in D / 2, t \mapsto \rho_{t}\left(z_{1}, z_{2}\right)$ is non-decreasing.
5) $\kappa_{d t}\left(z_{1}, z_{2}\right) \ll \rho_{d t}^{S}$, for every $z \in D$.
6) $\rho_{d t}\left(z_{1}, z_{2}, y_{1}, y_{2}\right) \ll \rho_{d t}^{S}$, for every $\left(z_{1}, z_{2}\right),\left(y_{1}, y_{2}\right) \in D / 2$.

Remark 2.25. Note that, for any $\left(z_{1}, z_{2}\right) \in D, X^{z_{1}} S^{z_{2}}$ is a special semimartingale. Indeed, by Proposition 2.24, $X_{t}^{z_{1}} S_{t}^{z_{2}}=N_{t} e^{\kappa_{t}\left(z_{1}, z_{2}\right)}$ where $\kappa\left(z_{1}, z_{2}\right)$ is a bounded variation continuous function and $N$ is a martingale. Hence, integration by parts implies that $X^{z_{1}} S^{z_{2}}$ is a special semimartingale whose decomposition is given by

$$
\begin{equation*}
X^{z_{1}} S^{z_{2}}=M\left(z_{1}, z_{2}\right)+V\left(z_{1}, z_{2}\right) \tag{2.31}
\end{equation*}
$$

where $M_{t}\left(z_{1}, z_{2}\right)=X_{0}^{z_{1}} S_{0}^{z_{2}}+\int_{0}^{t} e^{\kappa_{u}\left(z_{1}, z_{2}\right)} d N_{u}$ and $V_{t}\left(z_{1}, z_{2}\right)=\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}} \kappa_{d u}\left(z_{1}, z_{2}\right), \forall t \in$ $[0, T]$.

The following proposition shows that the local martingale part of the decomposition above is a square integrable martingale if $\left(z_{1}, z_{2}\right) \in D / 2$ and gives its angle bracket in terms of the generating function.

Proposition 2.26. Let $z=\left(z_{1}, z_{2}\right), y=\left(y_{1}, y_{2}\right) \in D / 2$. Then $X^{z_{1}} S^{z_{2}}$ is a special semimartingale, whose decomposition $X^{z_{1}} S^{z_{2}}=M\left(z_{1}, z_{2}\right)+V\left(z_{1}, z_{2}\right)$ satisfies, for $t \in[0, T]$,

$$
\begin{aligned}
V\left(z_{1}, z_{2}\right)_{t} & =\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}} \kappa_{d u}\left(z_{1}, z_{2}\right) \\
\left\langle M\left(z_{1}, z_{2}\right), M\left(y_{1}, y_{2}\right)\right\rangle_{t} & =\int_{0}^{t} X_{u-}^{z_{1}+y_{1}} S_{u-}^{z_{2}+y_{2}} \rho_{d u}\left(z_{1}, z_{2}, y_{1}, y_{2}\right) .
\end{aligned}
$$

In particular,

$$
\left\langle M\left(z_{1}, z_{2}\right)\right\rangle_{t}:=\left\langle M\left(z_{1}, z_{2}\right), \overline{M\left(z_{1}, z_{2}\right)}\right\rangle_{t}=\int_{0}^{t} X_{u-}^{2 \operatorname{Re}\left(z_{1}\right)} S_{u-}^{2 \operatorname{Re}\left(z_{2}\right)} \rho_{d u}\left(z_{1}, z_{2}\right)
$$

Moreover, $M\left(z_{1}, z_{2}\right)$ is a square integrable martingale.
Proof. This can be done adapting the techniques of [19, Lemma 3.2] and its generalization to one-dimensional additive processes, i.e. [17, Proposition 3.17 and Lemma 13.19].

The measure $d \rho^{S}$, called reference variance measure in [17], plays a central role in the expression of the canonical decomposition of special semimartingales depending on the couple $(X, S)$.

Corollary 2.27. The semimartingale decomposition of $S$ is given by $S=M^{S}+V^{S}$, where, for $t \in[0, T]$

$$
\begin{align*}
V_{t}^{S} & =\int_{0}^{t} S_{u-} \kappa_{d u}(0,1)  \tag{2.32}\\
\left\langle M^{S}\right\rangle_{t} & =\int_{0}^{t} S_{u-}^{2} \rho_{d u}^{S} \tag{2.33}
\end{align*}
$$

Proof. It follows from Proposition 2.26 setting $z_{1}=0, z_{2}=1$.
Now we state some useful estimates.
Lemma 2.28. Let $(a, b) \in D \cap \mathbb{R}^{2}$. Then

$$
\mathbb{E}\left[\sup _{t \leq T} X_{t}^{a} S_{t}^{b}\right]<\infty
$$

Proof. Let $(a, b) \in D \cap \mathbb{R}^{2}$, then $(a / 2, b / 2) \in D / 2$. By Proposition 2.26 , we have

$$
X_{t}^{a / 2} S_{t}^{b / 2}=M_{t}(a / 2, b / 2)+\int_{0}^{t} X_{u-}^{a / 2} S_{u-}^{b / 2} \kappa_{d u}(a / 2, b / 2), \forall t \in[0, T]
$$

and $M(a / 2, b / 2)$ is a square integrable martingale. Hence, by Doob inequality, we have

$$
\mathbb{E}\left[\sup _{t \leq T}\left|M_{t}(a / 2, b / 2)\right|^{2}\right] \leq 4 \mathbb{E}\left[\left|M_{T}(a / 2, b / 2)\right|^{2}\right]<\infty
$$

On the other hand, using Cauchy-Schwarz inequality and Fubini theorem, we obtain

$$
\begin{aligned}
\mathbb{E}\left[\sup _{t \leq T}\left|\int_{0}^{t} X_{u-}^{a / 2} S_{u-}^{b / 2} \kappa_{d u}(a / 2, b / 2)\right|^{2}\right] & \leq\|\kappa(a / 2, b / 2)\|_{T} \int_{0}^{T} \mathbb{E}\left[X_{u-}^{a} S_{u-}^{b}\right]\|\kappa(a / 2, b / 2)\|_{d u} \\
& =\|\kappa(a / 2, b / 2)\|_{T} \int_{0}^{T} e^{\kappa_{u}(a, b)}\|\kappa(a / 2, b / 2)\|_{d u} \\
& \leq e^{\|\kappa(a, b)\|_{T}}\|\kappa(a / 2, b / 2)\|_{T}^{2}<\infty
\end{aligned}
$$

Finally

$$
\mathbb{E}\left[\sup _{t \leq T} X_{t}^{a} S_{t}^{b}\right]=\mathbb{E}\left[\sup _{t \leq T}\left|X_{t}^{a / 2} S_{t}^{b / 2}\right|^{2}\right]<\infty
$$

In the general case, when $\left(z_{1}, z_{2}\right) \in D$, the local martingale part of the special semimartingale $X^{z_{1}} S^{z_{2}}$ is a true (not necessarily square integrable) martingale.

Proposition 2.29. Let $\left(z_{1}, z_{2}\right) \in D$, then, $M\left(z_{1}, z_{2}\right)$, the local martingale part of $X^{z_{1}} S^{z_{2}}$, is a true martingale such that

$$
\mathbb{E}\left[\sup _{t \leq T}\left|M_{t}\left(z_{1}, z_{2}\right)\right|\right]<\infty
$$

Proof. Let $\left(z_{1}, z_{2}\right) \in D$. Adopting the notations of (2.31), we remind that, by Proposition 2.26, $\forall t \in[0, T], M_{t}\left(z_{1}, z_{2}\right)=X_{t}^{z_{1}} S_{t}^{z_{2}}-\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}} \kappa_{d u}\left(z_{1}, z_{2}\right)$. For this local martingale we can write

$$
\begin{aligned}
\mathbb{E}\left[\sup _{t \leq T}\left|M_{t}\left(z_{1}, z_{2}\right)\right|\right] & \leq \mathbb{E}\left[\sup _{t \leq T}\left|X_{t}^{z_{1}} S_{t}^{z_{2}}\right|\right]+\mathbb{E}\left[\int_{0}^{T}\left|X_{t-}^{z_{1}} S_{t-}^{z_{2}}\right|\left\|\kappa\left(z_{1}, z_{2}\right)\right\|_{d t}\right] \\
& \leq \mathbb{E}\left[\sup _{t \leq T}\left|X_{t}^{\operatorname{Re}\left(z_{1}\right)} S_{t}^{\operatorname{Re}\left(z_{2}\right) \mid}\right|\right]\left(1+\left\|\kappa\left(z_{1}, z_{2}\right)\right\|_{T}\right) .
\end{aligned}
$$

Since $\left(\operatorname{Re}\left(z_{1}\right), \operatorname{Re}\left(z_{2}\right)\right)$ belongs to $D$, by Lemma 2.28 , the right-hand side is finite. Consequently the local martingale $M\left(z_{1}, z_{2}\right)$ is indeed a true martingale.

The goal of this section is to show that $(X, S)$ is a solution of a strong martingale problem, with related triplet $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$, which will be specified below. For this purpose, we determine the semimartingale decomposition of $\left(f\left(t, X_{t}, S_{t}\right)\right)_{t}$ for functions $f:[0, T] \times \mathcal{O} \rightarrow \mathbb{C}$, where $\mathcal{O}=] 0, \infty\left[^{2}\right.$, of the form

$$
\begin{equation*}
f(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right), \forall t \in[0, T], x, y>0 \tag{2.34}
\end{equation*}
$$

where $\Pi$ is finite complex Borel measure on $\mathbb{C}^{2}$ and $\lambda:[0, T] \times \mathbb{C}^{2} \longrightarrow \mathbb{C}$. The family of those functions will include the set $\mathfrak{D}(\mathfrak{a})$ defined later.

Proposition 2.29 and item 5) of Proposition 2.24 say that, for $z=\left(z_{1}, z_{2}\right) \in D$,

$$
t \mapsto X_{t}^{z_{1}} S_{t}^{z_{2}}-\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}} \kappa_{d u}\left(z_{1}, z_{2}\right)=X_{t}^{z_{1}} S_{t}^{z_{2}}-\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}} \frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}} \rho_{d u}^{S}
$$

is a martingale. This provides the semimartingale decomposition of the basic functions $(t, x, s) \mapsto x^{z_{1}} z^{z_{2}}$ for $z_{1}, z_{2} \in D$, applied to $(X, S)$. Those functions are expected to be elements of $\mathfrak{D}(\mathfrak{a})$ and one candidate for the bounded variation process $A$ is $\rho^{S}$. It remains to precisely define $\mathfrak{D}(\mathfrak{a})$ and the operator $\mathfrak{a}$.

A first step in this direction is to consider a Borel function $\lambda:[0, T] \times \mathbb{C}^{2} \rightarrow \mathbb{C}$ such that, for any $\left(z_{1}, z_{2}\right) \in D, t \in[0, T] \mapsto \lambda\left(t, z_{1}, z_{2}\right)$ is absolutely continuous with respect to $\rho^{S}$.

Lemma 2.30. Let $\lambda:[0, T] \times \mathbb{C}^{2} \rightarrow \mathbb{C}$ such that, for any $\left(z_{1}, z_{2}\right) \in D, t \in[0, T] \mapsto \lambda\left(t, z_{1}, z_{2}\right)$ is absolutely continuous with respect to $\rho^{S}$. Then for any $\left(z_{1}, z_{2}\right) \in D$,
$t \mapsto M_{t}^{\lambda}\left(z_{1}, z_{2}\right):=S_{t}^{z_{1}} X_{t}^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right)-\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}}\left\{\frac{d \lambda\left(u, z_{1}, z_{2}\right)}{d \rho_{u}^{S}}+\lambda\left(u, z_{1}, z_{2}\right) \frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}\right\} \rho_{d u}^{S}$,
is a martingale. Moreover, if $\left(z_{1}, z_{2}\right) \in D / 2$ then $M^{\lambda}\left(z_{1}, z_{2}\right)$ is a square integrable martingale and

$$
\begin{equation*}
\mathbb{E}\left[\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|^{2}\right]=\int_{0}^{t} e^{\kappa_{u}\left(2 \operatorname{Re}\left(z_{1}\right), 2 \operatorname{Re}\left(z_{2}\right)\right)}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right) \tag{2.36}
\end{equation*}
$$

Proof. Let $\left(z_{1}, z_{2}\right) \in D, M\left(z_{1}, z_{2}\right)$ and $V\left(z_{1}, z_{2}\right)$ be the random fields introduced in Remark 2.25. Since $\lambda\left(d t, z_{1}, z_{2}\right) \ll \rho_{d t}^{S}$, then $t \mapsto \lambda\left(t, z_{1}, z_{2}\right)$ is a bounded continuous function on
$[0, T]$. By item 5) of Proposition $2.24 M^{\lambda}\left(z_{1}, z_{2}\right)$ is well-defined. Integrating by parts and taking into account Remark 2.25 allows to show

$$
\begin{equation*}
M_{t}^{\lambda}\left(z_{1}, z_{2}\right)=\lambda\left(0, z_{1}, z_{2}\right) M_{0}\left(z_{1}, z_{2}\right)+\int_{0}^{t} \lambda\left(u, z_{1}, z_{2}\right) d M_{u}\left(z_{1}, z_{2}\right), \forall t \in[0, T] \tag{2.37}
\end{equation*}
$$

Obviously $M^{\lambda}\left(z_{1}, z_{2}\right)$ is a local martingale. In order to prove that it is a true martingale, we establish that

$$
\mathbb{E}\left[\sup _{t \leq T}\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|\right]<\infty
$$

Indeed, by integration by parts in (2.37), for $t \in[0, T]$ we have

$$
M_{t}^{\lambda}\left(z_{1}, z_{2}\right)=\lambda\left(t, z_{1}, z_{2}\right) M_{t}\left(z_{1}, z_{2}\right)-\int_{0}^{t} M_{u-}\left(z_{1}, z_{2}\right) \lambda\left(d u, z_{1}, z_{2}\right)
$$

Hence, as in the proof of Lemma 2.28,

$$
\begin{align*}
\mathbb{E}\left[\sup _{t \leq T}\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|\right] & \leq \mathbb{E}\left[\sup _{t \leq T}\left|\lambda\left(t, z_{1}, z_{2}\right) M_{t}\left(z_{1}, z_{2}\right)\right|\right]+\mathbb{E}\left[\int_{0}^{T}\left|M_{u-}\left(z_{1}, z_{2}\right)\right|\left\|\lambda\left(., z_{1}, z_{2}\right)\right\|_{d t}\right] \\
& \leq 2 \mathbb{E}\left[\sup _{t \leq T}\left|M_{t}\left(z_{1}, z_{2}\right)\right|\right]\left\|\lambda\left(., z_{1}, z_{2}\right)\right\|_{T} \tag{2.38}
\end{align*}
$$

Thanks to Proposition 2.29 , the right-hand side of $(2.38)$ is finite and finally $M^{\lambda}\left(z_{1}, z_{2}\right)$ is shown to be a martingale so that the first part of Lemma 2.30 is proved.

Now, suppose that $\left(z_{1}, z_{2}\right) \in D / 2$. By (2.37) and Proposition 2.26, we have

$$
\begin{align*}
\mathbb{E}\left[\left\langle M^{\lambda}\left(z_{1}, z_{2}\right)\right\rangle_{T}\right] & =\mathbb{E}\left[\int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2}\left\langle M\left(z_{1}, z_{2}\right)\right\rangle_{d t}\right] \\
& =\mathbb{E}\left[\int_{0}^{T} X_{u-}^{2 \operatorname{Re}\left(z_{1}\right)} S_{u-}^{2 \operatorname{Re}\left(z_{2}\right)}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right)\right] \\
& =\int_{0}^{T} e^{\kappa_{u}\left(2 \operatorname{Re}\left(z_{1}\right), 2 \operatorname{Re}\left(z_{2}\right)\right)}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right)  \tag{2.39}\\
& \leq \sup _{u \leq T}^{\kappa^{\kappa}\left(2 \operatorname{Re}\left(z_{1}\right), 2 \operatorname{Re}\left(z_{2}\right)\right)} \int_{0}^{T}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right)<\infty .
\end{align*}
$$

The latter term is finite by point 3 ) of Proposition 2.24 and by the fact that $\lambda\left(., z_{1}, z_{2}\right)$ is bounded on $[0, T]$. Consequently, $M^{\lambda}\left(z_{1}, z_{2}\right)$ is a square integrable martingale and since $\left|M^{\lambda}\left(z_{1}, z_{2}\right)\right|^{2}-\left\langle M^{\lambda}\left(z_{1}, z_{2}\right)\right\rangle$ is a martingale, then

$$
\mathbb{E}\left[\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|^{2}\right]=\int_{0}^{t} e^{\kappa_{u}\left(2 \operatorname{Re}\left(z_{1}\right), 2 \operatorname{Re}\left(z_{2}\right)\right)}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right)
$$

because of (2.39).
Now, let $\Pi$ be a finite Borel measure on $\mathbb{C}^{2}$ and let us make the following assumption on it.

Assumption 2.31. We set $I_{0}:=\operatorname{Re}(\operatorname{supp} \Pi)$.

1. $I_{0}$ is bounded.
2. $I_{0} \subset D$.

Note that this assumption implies that supp $\Pi \subset D$.
Theorem 2.32. Suppose that Assumptions 2.23 and 2.31 are verified. Let $\lambda:[0, T] \times \mathbb{C}^{2} \rightarrow \mathbb{C}$ be a function such that

$$
\begin{align*}
\forall\left(z_{1}, z_{2}\right) \in \operatorname{supp} \Pi, \lambda\left(d t, z_{1}, z_{2}\right) & \ll \rho_{d t}^{S}  \tag{2.40}\\
\forall t \in[0, T], \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} & <\infty  \tag{2.41}\\
\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right| & <\infty \tag{2.42}
\end{align*}
$$

Then the function $f$ defined by

$$
\begin{equation*}
f(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right), \forall t \in[0, T], x, s>0 \tag{2.43}
\end{equation*}
$$

is continuous. Moreover
$t \mapsto M_{t}^{\lambda}:=f\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \rho_{d u}^{S} \int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) X_{u-}^{z_{1}} S_{u-}^{z_{2}}\left\{\frac{d \lambda\left(u, z_{1}, z_{2}\right)}{d \rho_{u}^{S}}+\lambda\left(u, z_{1}, z_{2}\right) \frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}\right\}$
is a martingale.

Remark 2.33. In (2.42) and (2.44), part of the integrand with respect to $\Pi$ is only defined for $\left(z_{1}, z_{2}\right) \in D$. By convention the integrand will be set to zero for $\left(z_{1}, z_{2}\right) \notin D$. In the sequel we will adopt the same rule.

Proof. Let $\lambda:[0, T] \times \mathbb{C}^{2} \rightarrow \mathbb{C}$ verifying the hypotheses of the theorem.
The function $f$ is well-defined. Indeed, for $t \in[0, T], x, y>0$,

$$
|f(t, x, s)| \leq \sup _{(a, b) \in I_{0}} x^{a} s^{b} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right)\right|
$$

which is finite because of Condition (2.41) and Assumption 2.31, taking into account CauchySchwarz inequality.

Moreover, by Fubini theorem and (2.43), we get

$$
\begin{align*}
\mathbb{E}\left[\left|f\left(t, X_{t}, S_{t}\right)\right|\right] & \leq \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \mathbb{E}\left[X_{t}^{\operatorname{Re}\left(z_{1}\right)} S_{t}^{\operatorname{Re}\left(z_{2}\right)}\right]\left|\lambda\left(t, z_{1}, z_{2}\right)\right| \\
& \leq \sup _{u \in[0, T],(a, b) \in I_{0}} \mathbb{E}\left[X_{u}^{a} S_{u}^{b}\right] \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right)\right| \tag{2.45}
\end{align*}
$$

The right-hand side is finite by item 3) of Proposition 2.24 and Condition (2.41).

We observe that $t \mapsto \lambda\left(t, z_{1}, z_{2}\right)$ is a continuous function since it is absolutely continuous with respect to $\rho^{S}$ for fixed $\left(z_{1}, z_{2}\right) \in \mathbb{C}^{2}$. On the other hand, Condition (2.41) implies that the family $\left(\lambda\left(t, z_{1}, z_{2}\right), t \in[0, T]\right)$ is $|\Pi|$-uniformly integrable. These properties, together with Lebesgue dominated convergence theorem imply that $f$ defined in (2.43) is continuous with respect to all the variables.

We show now that the process $t \mapsto M_{t}^{\lambda}$ is well-defined. This holds because

$$
\begin{align*}
& \mathbb{E}\left[\int_{0}^{t} \rho_{d u}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|X_{u-}^{z_{1}} S_{u-}^{z_{2}}\right|\left|\frac{d \lambda\left(u, z_{1}, z_{2}\right)}{d \rho_{u}^{S}}+\lambda\left(u, z_{1}, z_{2}\right) \frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}\right|\right]  \tag{2.46}\\
\leq & \sup _{u \in[0, T],(a, b) \in I_{0}} \mathbb{E}\left[X_{u}^{a} S_{u}^{b}\right] \int_{0}^{t} \rho_{d u}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(u, z_{1}, z_{2}\right)}{d \rho_{u}^{S}}+\lambda\left(u, z_{1}, z_{2}\right) \frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}\right|,
\end{align*}
$$

which is finite by point 3 ) of Proposition 2.24 and Condition (2.42).
(2.46) allows to apply Fubini theorem to the integral term in (2.44), so that we get

$$
\begin{align*}
M_{t}^{\lambda} & =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right)\left(X_{t}^{z_{1}} S_{t}^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right)-\int_{0}^{t} X_{u-}^{z_{1}} S_{u-}^{z_{2}}\left\{\frac{d \lambda\left(u, z_{1}, z_{2}\right)}{d \rho_{u}^{S}}+\lambda\left(u, z_{1}, z_{2}\right) \frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}\right\} \rho_{d u}^{S}\right) \\
& =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) M_{t}^{\lambda}\left(z_{1}, z_{2}\right), \tag{2.47}
\end{align*}
$$

where $M^{\lambda}\left(z_{1}, z_{2}\right)$ is defined in (2.35) for any $\left(z_{1}, z_{2}\right) \in D$. We observe that

$$
\begin{equation*}
\mathbb{E}\left[\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|\right]<\infty \tag{2.48}
\end{equation*}
$$

taking into account (2.45) and (2.46). It remains to show that $M^{\lambda}$ is a martingale.
Let $0 \leq s \leq t \leq T$ and a bounded, $\mathcal{F}_{s}$-measurable random variable $G$. By Fubini theorem and Lemma 2.30 it follows

$$
\begin{aligned}
\mathbb{E}\left[M_{t}^{\lambda} G\right] & =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) \mathbb{E}\left[M_{t}^{\lambda}\left(z_{1}, z_{2}\right) G\right] \\
& =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) \mathbb{E}\left[M_{s}^{\lambda}\left(z_{1}, z_{2}\right) G\right] \\
& =\mathbb{E}\left[\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) M_{s}^{\lambda}\left(z_{1}, z_{2}\right) G\right] \\
& =\mathbb{E}\left[M_{s}^{\lambda} G\right]
\end{aligned}
$$

which implies the desired result.
We proceed now to the definition of the domain $\mathfrak{D}(\mathfrak{a})$ in view of the specification of the corresponding martingale problem. We set
$\mathfrak{D}(\mathfrak{a})=\left\{f:(t, x, s) \mapsto \int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right), \forall t \in[0, T], x, y>0\right.$,
where $\Pi$ is a finite Borel measure on $\mathbb{C}^{2}$ verifying Assumption 2.31, with $\lambda:[0, T] \times \mathbb{C}^{2} \rightarrow \mathbb{C}$ Borel such that conditions (2.40), (2.41) and (2.42) are fulfilled $\}$.

Corollary 2.34. Suppose that Assumptions 2.23 and 2.31 are verified. Then $(X, S)$ is a solution of the strong martingale problem related to $\left(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, \rho^{S}\right)$ where, for $f \in \mathfrak{D}(\mathfrak{a})$ of the type (2.43), $\mathfrak{a}(f)$ is defined by

$$
\begin{equation*}
\mathfrak{a}(f)(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}}\left\{\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right\} \tag{2.50}
\end{equation*}
$$

for all $t \in[0, T], x, s>0$.
Proof. By Theorem 2.32 note that $f \in \mathfrak{D}(\mathfrak{a})$ defined in (2.43) is continuous, which implies that (2.2) is fulfilled. By (2.42), $\mathfrak{a}(f)$ belongs to $\mathcal{L}$ defined in (2.3) and Condition (2.4) is fulfilled. Finally (2.5) is a consequence of (2.44) in Theorem 2.32.

Under additional conditions, one can say more about the martingale decomposition given by the strong martingale problem related to $\left(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, \rho^{S}\right)$.

Proposition 2.35. Let $f \in \mathfrak{D}(\mathfrak{a})$ as defined in (2.43). Suppose the following.
a) $I_{0}:=\operatorname{Re}(\operatorname{supp} \Pi) \subset D / 2$,
b) $\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \int_{0}^{T}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right)<\infty$.

Then, the martingale $t \mapsto M_{t}^{\lambda}=f\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(f)\left(u, X_{u-}, S_{u-}\right) \rho_{d u}^{S}$ is square integrable.
Proof. Let $t \in[0, T]$ and $M^{\lambda}$ as defined in (2.44), which is a martingale by Theorem 2.32. By (2.47) we have

$$
\begin{equation*}
M_{t}^{\lambda}=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) M_{t}^{\lambda}\left(z_{1}, z_{2}\right) \tag{2.51}
\end{equation*}
$$

where $M^{\lambda}\left(z_{1}, z_{2}\right)$ was defined in (2.35). By Lemma 2.30 , for every $\left(z_{1}, z_{2}\right) \in D / 2$, we have

$$
\begin{equation*}
\mathbb{E}\left[\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|^{2}\right]=\int_{0}^{t} e^{\kappa_{u}\left(2 \operatorname{Re}\left(z_{1}\right), 2 \operatorname{Re}\left(z_{2}\right)\right)}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right) \tag{2.52}
\end{equation*}
$$

By Fubini theorem, integrating both sides of (2.52) with respect to $|\Pi|$, gives

$$
\begin{aligned}
\mathbb{E}\left[\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|^{2}\right] & =\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \mathbb{E}\left[\left|M_{t}^{\lambda}\left(z_{1}, z_{2}\right)\right|^{2}\right] \\
& =\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \int_{0}^{t} e^{\kappa_{u}\left(2 \operatorname{Re}\left(z_{1}\right), 2 \operatorname{Re}\left(z_{2}\right)\right)}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right) \\
& \leq \sup _{u \in[0, T],(a, b) \in I_{0}} e^{\kappa_{u}(a, b)} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \int_{0}^{t}\left|\lambda\left(u, z_{1}, z_{2}\right)\right|^{2} \rho_{d u}\left(z_{1}, z_{2}\right)
\end{aligned}
$$

Now, by point 3) of Proposition 2.24 and condition b), the right-hand side is finite. This together with (2.51) and Cauchy-Schwarz show that $M^{\lambda}$ is square integrable.

Proposition 2.36. We suppose the validity of Assumptions 2.23.

1) Assumption 2.4 is verified. More precisely
i) id: $(t, x, s) \longmapsto s \in \mathfrak{D}(\mathfrak{a})$ and

$$
\begin{equation*}
\mathfrak{a}(i d)(t, x, s)=s \frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}, \quad \forall t \in[0, T], x, s>0 \tag{2.53}
\end{equation*}
$$

ii) $(t, x, s) \longmapsto s^{2} \in \mathfrak{D}(\mathfrak{a})$ and

$$
\begin{equation*}
\tilde{\mathfrak{a}}(i d)(t, x, s)=s^{2}, \quad \forall t \in[0, T], x, s>0 \tag{2.54}
\end{equation*}
$$

2) Let $\Pi$ be a finite signed Borel measure on $\mathbb{C}^{2}$ verifying Assumption 2.31. Let $f \in \mathfrak{D}(\mathfrak{a})$ of the form (2.49), such that $\widetilde{f}=f \times i d \in \mathfrak{D}(\mathfrak{a})$. Then $\tilde{\mathfrak{a}}$, defined in (2.7), is given by, $\forall t \in[0, T], x, s>0$,

$$
\begin{equation*}
\widetilde{\mathfrak{a}}(f)(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) \lambda\left(t, z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}+1} \frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}} \tag{2.55}
\end{equation*}
$$

Proof.
We first address item 1).
i) Let $\Pi_{1}\left(z_{1}, z_{2}\right)=\delta_{\left\{z_{1}=0, z_{2}=1\right\}}$ and $\lambda \equiv 1$. Since by Assumption $2.23(0,1) \in D, \Pi_{1}$ fulfills Assumption 2.31. The other conditions (2.40), (2.41), (2.42) defining $\mathfrak{D}(\mathfrak{a})$ in (2.49) are trivially satisfied. Consequently $i d \in \mathfrak{D}(\mathfrak{a})$ and (2.53) follows from (2.50).
ii) Let $\Pi_{2}\left(z_{1}, z_{2}\right)=\delta_{\left\{z_{1}=0, z_{2}=2\right\}}$ and $\lambda \equiv 1$. Again, by Assumption $2.23(0,2) \in D$, and $\Pi_{2}$ fulfills Assumption 2.31. Similar arguments as for i) allow to show that $(t, x, s) \longmapsto$ $s^{2} \in \mathfrak{D}(\mathfrak{a})$.

Formula (2.55) constitutes a direct application of (2.50), taking into account (2.49), which establishes item 2). In particular (2.54) holds.

## 3 The basic BSDE and the deterministic problem

### 3.1 General framework

We consider two $\mathcal{F}_{t}$-adapted processes $(X, S)$ fulfilling the martingale problem related to $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$ stated in Definition 2.1 under Assumption 2.4. We denote by $M^{S}$ the martingale part of $S$.

Let $\hat{f}: \Omega \times[0, T] \times \mathbb{C}^{2} \longrightarrow \mathbb{C}$ be a predictable random field (i.e. such that for every $y, z$, $s \mapsto \hat{f}(\cdot, s, y, z)$ is predictable) and $h$ be an $\mathcal{F}_{T}$-measurable, complex valued, random variable. As we have mentioned in the introduction, the object of our interest is a BSDE of the type (1.2). We focus on a deterministic natural problem associated with it, which plays the role of the semilinear PDE of the Brownian case.

Definition 3.1. A triplet $(Y, Z, O)$ of processes is called solution of (1.2) if the following holds.

1) $\left(Y_{t}\right)$ is $\mathcal{F}_{t^{-}}$-adapted;
2) $\left(Z_{t}\right)$ is $\mathcal{F}_{t}$-predictable and
(a) $\int_{0}^{T}\left|Z_{s}\right|^{2} d\left\langle M^{S}\right\rangle_{s}<\infty$ a.s.
(b) $\int_{0}^{T}\left|\hat{f}\left(\omega, s, Y_{s-}, Z_{s}\right)\right| d\left\|V^{S}\right\|_{s}<\infty$ a.s.
3) Equality (1.2) holds and $\left(O_{t}\right)$ is an $\mathcal{F}_{t}$-local martingale such that $\left\langle O, M^{S}\right\rangle=0$ and $O_{0}=0$ a.s.

In this section we are more particularly interested in the $\operatorname{BSDE}(1.2)$ when $\hat{f}$ is given by (1.3).

### 3.2 The forward-backward case and the deterministic problem

As we have already mentioned in the Introduction, the BSDE on which we will focus on, arises when the driver coefficient $\hat{f}$ is associated with a locally bounded function $f:[0, T] \times \mathcal{O} \times$ $\mathbb{C}^{2} \longrightarrow \mathbb{C}$ and with the $\mathcal{F}_{t}$-special semimartingale $(X, S)$ which solves the strong martingale problem related to $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$. In conformity with (1.3), we suppose the form of $\hat{f}$ and of the target r.v. $h$ as follows.

$$
\begin{aligned}
\mathfrak{a}(i d)\left(t, X_{t-}(\omega), S_{t-}(\omega)\right) \hat{f}(\omega, t, y, z) & =f\left(t, X_{t-}(\omega), S_{t-}(\omega), y, z\right) \\
h & =g\left(X_{T}, S_{T}\right)
\end{aligned}
$$

for some continuous function $g: \mathcal{O} \rightarrow \mathbb{C}$.
Therefore, our subject of study is the particular BSDE given below.

$$
\begin{equation*}
Y_{t}=h+\int_{t}^{T} f\left(r, X_{r-}, S_{r-}, Y_{r-}, Z_{r}\right) d A_{r}-\int_{t}^{T} Z_{r} d M_{r}^{S}-\left(O_{T}-O_{t}\right), \quad t \in[0, T] \tag{3.1}
\end{equation*}
$$

As we remarked in the introduction, when $M^{S}$ is a Brownian motion and $\left(\mathcal{F}_{t}\right)$ is its canonical filtration, (3.1) can be linked to a semilinear partial differential equation. We will formulate a deterministic problem, generalizing that "classical" semilinear PDE. In particular we look for solutions $(Y, Z, O)$ for which there is a function $y \in \mathfrak{D}(\mathfrak{a})$ such that $\widetilde{y}=y \times i d \in$ $\mathfrak{D}(\mathfrak{a})$ and a locally bounded Borel function $z:[0, T] \times \mathcal{O} \longrightarrow \mathbb{C}$, such that

$$
\begin{align*}
Y_{t} & =y\left(t, X_{t}, S_{t}\right)  \tag{3.2}\\
Z_{t} & =z\left(t, X_{t-}, S_{t-}\right), \quad \forall t \in[0, T] \tag{3.3}
\end{align*}
$$

and

$$
\begin{align*}
\int_{0}^{t}\left|Z_{s}\right|^{2} d\left\langle M^{S}\right\rangle_{s} & <\infty \text { a.s. } \\
\int_{0}^{t}\left|f\left(s, X_{s-}, S_{s-}, Y_{s-}, Z_{s}\right)\right| d\|A\|_{s} & <\infty \text { a.s. } \tag{3.4}
\end{align*}
$$

By (3.2) and (3.4), Conditions 1) and 2) of Definition 3.1 are obviously fulfilled. Consequently the triplet $(Y, Z, O)$ where

$$
\begin{equation*}
O_{t}:=Y_{t}-Y_{0}-\int_{0}^{t} Z_{r} d M_{r}^{S}+\int_{0}^{t} f\left(r, X_{r-}, S_{r-}, Y_{r-}, Z_{r}\right) d A_{r} \tag{3.5}
\end{equation*}
$$

is a solution of (1.2) provided that

$$
\begin{equation*}
\text { 1. }\left(O_{t}\right) \text { is an } \mathcal{F}_{t} \text {-local martingale, } \tag{3.6}
\end{equation*}
$$

2. $\left\langle O, M^{S}\right\rangle=0$,
3. $Y_{T}=g\left(X_{T}, S_{T}\right)$.

Since $(X, S)$ solves the strong martingale problem related to ( $\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$, replacing (3.2) in expression (3.5), Condition (3.6) can be rewritten saying that

$$
\int_{0}^{t} \mathfrak{a}(y)\left(r, X_{r-}, S_{r-}\right) d A_{r}+\int_{0}^{t} f\left(r, X_{r-}, S_{r-}, Y_{r-}, Z_{r}\right) d A_{r}
$$

is an $\mathcal{F}_{t}$-local martingale. This implies that

$$
\begin{equation*}
\int_{0}^{t} \mathfrak{a}(y)\left(r, X_{r-}, S_{r-}\right) d A_{r}+\int_{0}^{t} f\left(r, X_{r-}, S_{r-}, Y_{r-}, Z_{r}\right) d A_{r}=0 \tag{3.9}
\end{equation*}
$$

On the other hand, Condition (3.7) implies

$$
\begin{equation*}
\left\langle M^{Y}, M^{S}\right\rangle_{t}=\int_{0}^{t} Z_{s} d\left\langle M^{S}\right\rangle_{s} \tag{3.10}
\end{equation*}
$$

where $M^{Y}$ denotes the martingale part of $Y$. By Lemma 2.3 and item ii) of Corollary 2.5, we have

$$
\begin{aligned}
\left\langle M^{Y}, M^{S}\right\rangle_{t} & =\int_{0}^{t} \widetilde{\mathfrak{a}}(y)\left(r, X_{r-}, S_{r-}\right) d A_{r} \\
\left\langle M^{S}\right\rangle_{t} & =\int_{0}^{t} \widetilde{\mathfrak{a}}(i d)\left(r, X_{r-}, S_{r-}\right) d A_{r}
\end{aligned}
$$

Consequently, Condition (3.10) can be re-expressed as

$$
\begin{equation*}
\int_{0}^{t} \widetilde{\mathfrak{a}}(y)\left(r, X_{r-}, S_{r-}\right) d A_{r}=\int_{0}^{t} z\left(r, X_{r-}, S_{r-}\right) \widetilde{\mathfrak{a}}(i d)\left(r, X_{r-}, S_{r-}\right) d A_{r} \tag{3.11}
\end{equation*}
$$

Condition (3.8) requires $y(T, \cdot, \cdot)=g$.
This allows to state the following representation theorem.
Theorem 3.2. Suppose the existence of a function $y$, such that $y, \widetilde{y}:=y \times i d$ belong to $\mathfrak{D}(\mathfrak{a})$, and a Borel locally bounded function $z$, solving the system

$$
\begin{align*}
\mathfrak{a}(y)(t, x, s) & =-f(t, x, s, y(t, x, s), z(t, x, s))  \tag{3.12}\\
\widetilde{\mathfrak{a}}(y)(t, x, s) & =z(t, x, s) \widetilde{\mathfrak{a}}(i d)(t, x, s) \tag{3.13}
\end{align*}
$$

for $t \in[0, T]$ and $(x, s) \in \mathcal{O}$, where the equalities hold in $\mathcal{L}$, with the terminal condition $y(T, .,)=.g(.,$.$) .$

Then the triplet $(Y, Z, O)$ defined by

$$
\begin{equation*}
Y_{t}=y\left(t, X_{t}, S_{t}\right), Z_{t}=z\left(t, X_{t-}, S_{t-}\right) \tag{3.14}
\end{equation*}
$$

and $\left(O_{t}\right)$ given by (3.5), is a solution to the BSDE (3.1).
Proof. The triplet ( $Y, Z, O$ ) fulfills the three conditions of Definition 3.1 provided that (3.4) is verified. Indeed, since $y \in \mathfrak{D}(\mathfrak{a})$ then the integral $\int_{0}^{t}\left|f\left(s, X_{s-}, S_{s-}, Y_{s-}, Z_{s}\right)\right| d\|A\|_{s}$, is finite taking into account (2.4).
Since $z$ is locally bounded, then $\int_{0}^{T}\left|Z_{s}\right|^{2} d\langle M\rangle_{s}$ is also finite. This concludes the proof of the theorem.

## Remark 3.3.

1. The statement of Theorem 3.2 can be generalized relaxing the assumption on $z$ to be locally bounded. We replace this with the condition

$$
\begin{equation*}
\int_{0}^{T} z^{2}\left(r, X_{r-}, S_{r-}\right) \widetilde{\mathfrak{a}}(i d)\left(r, X_{r-}, S_{r-}\right) d A_{r}<\infty \text { a.s. } \tag{3.15}
\end{equation*}
$$

This is equivalent to $\int_{0}^{T}\left|Z_{s}\right|^{2} d\langle M\rangle_{s}<\infty$ a.s.
2. In particular, if $z$ is locally bounded a.s., then (3.15) is fulfilled.

Remark 3.4. Theorem 3.2 constitutes also an existence theorem for particular BSDEs. If $M^{S}$ is a square integrable martingale and the function $\hat{f}$ associated with $f$, fulfills some Lipschitz type conditions then the solution $(Y, Z, O)$ provided by (3.14) is unique in the class of processes introduced in [6, Theorem 3.1].

The presence of the local martingale $O$ is closely related to the classical martingale representation property. In fact, if $(\Omega, \mathcal{F}, \mathbb{P})$ verifies the local martingale representation property with respect to $M^{S}$, then $O$ vanishes.

Proposition 3.5. Suppose that $(\Omega, \mathcal{F}, \mathbb{P})$ fulfills the local martingale representation property with respect to $M$. Then, if $(Y, Z, O)$ is a solution to (3.1), then, necessarily $O_{t}=0, \forall t \in$ $[0, T]$.
Proof. Since $\left(O_{t}\right)$ is an $\mathcal{F}_{t}$-local martingale, there is a predictable process $\left(Z_{t}\right)$ such that

$$
O_{t}=O_{0}+\int_{0}^{t} Z_{s} d M_{s}^{S}, \forall t \in[0, T]
$$

So the condition $\left\langle O, M^{S}\right\rangle \equiv 0$ implies

$$
\int_{0} Z_{s} d\left\langle M^{S}\right\rangle=0 .
$$

Consequently,

$$
Z \equiv 0 d \mathbb{P} \otimes d\left\langle M^{S}\right\rangle \text { a.e., }
$$

and so $O_{t}=O_{0}=0 \forall t \in[0, T]$.

### 3.3 Illustration 1: the Markov semigroup case

Let us consider the case of Section 2.2 with related notations. Let $S=X^{0, x}$ be a solution of the strong martingale problem related to $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$, see Definition 2.6. Let $\left(P_{t}\right)$ be the semigroup introduced in (2.15), fulfilling Assumption 2.7 with generator $L$ defined in Definition 2.11. Let $f:[0, T] \times \mathbb{R} \times \mathbb{C} \longrightarrow \mathbb{C}$ be a locally bounded function and a continuous function $g: \mathbb{R} \longrightarrow \mathbb{C}$.
Here we have of course $S=M^{S}+V^{S}$ where $V^{S}=\int_{0} \mathfrak{a}(i d)\left(r, S_{r-}\right) d r$ and $i d(s) \equiv s$.
Theorem 3.2 gives the following.
Proposition 3.6. Suppose the existence of a function $y:[0, T] \times \mathbb{R} \rightarrow \mathbb{C}$ and a Borel locally bounded function $z:[0, T] \times \mathbb{R} \rightarrow \mathbb{C}$ verifying the following.
i) $t \mapsto y(t, \cdot)($ resp. $\tilde{y}(t, \cdot))$ takes value in $D(L)$ and it is continuous with respect to the graph norm.
ii) $t \mapsto y(t, \cdot)($ resp. $\tilde{y}(t, \cdot))$ is of class $C^{1}$ with values in $E$.
iii) For $(t, x) \in[0, T] \times \mathbb{R}$,

$$
\begin{aligned}
\partial_{t} y(t, x)+L y(t, \cdot)(x) & =-f(t, x, y(t, x), z(t, x)), \\
z(t, \cdot) \widetilde{L}(i d) & =\widetilde{L} y(t, \cdot), \\
y(T, .) & =g,
\end{aligned}
$$

where $\widetilde{L} \varphi=L \tilde{\varphi}-\varphi \operatorname{Lid}-i d L \varphi$.
Then the triplet $(Y, Z, O)$, where

$$
\begin{gathered}
Y_{t}:=y\left(t, S_{t}\right), Z_{t}:=z\left(t, S_{t-}\right), \\
O_{t}:=Y_{t}-Y_{0}-\int_{0}^{t} Z_{r} d M_{r}^{S}+\int_{0}^{t} \hat{f}\left(r, \omega, Y_{r-}, Z_{r}\right) d V_{r}^{S}, t \in[0, T],
\end{gathered}
$$

is a solution of the BSDE

$$
Y_{t}=g\left(S_{T}\right)+\int_{t}^{T} \hat{f}\left(r, \omega, Y_{r-}, Z_{r}\right) d V_{r}^{S}-\int_{t}^{T} Z_{r} d M_{r}^{S}-\left(O_{T}-O_{t}\right), \quad t \in[0, T],
$$

in the sense of Definition 3.1, where

$$
\mathfrak{a}(i d)\left(r, S_{r-}(\omega)\right) \hat{f}(r, \omega, y, z)=f\left(r, S_{r-}(\omega), y, z\right) .
$$

Remark 3.7. If $S=\sigma W$ with $\sigma>0$ and $\varphi:[0, T] \times \mathbb{R} \rightarrow \mathbb{C}$ is of class $C^{1,2}$ then $\mathfrak{a}(\varphi)=$ $\partial_{t} \varphi+\frac{\sigma^{2}}{2} \partial_{s s} \varphi$ and $\widetilde{\mathfrak{a}}(\varphi)=\sigma^{2} \partial_{s} \varphi=\widetilde{\mathfrak{a}}(i d) \partial_{s} \varphi$.
In the case where $L$ is a generic generator, the formal quotient $\frac{\tilde{\mathfrak{a}}(\varphi)}{\mathfrak{a}(i d)}$ can be considered as a sort of generalized derivative.

### 3.4 Illustration 2: the diffusion case

Consider the case of where ( $X, S$ ) is diffusion process as given in equations (2.28). We remind that in that case, the operator $\mathfrak{a}$, for $\varphi \in \mathcal{C}^{1,2}\left([0, T] \times \mathbb{R}^{2}\right)$, is given by

$$
\begin{aligned}
\mathfrak{a}(\varphi) & =\partial_{t} \varphi+b_{S} \partial_{s} \varphi+b_{X} \partial_{x} \varphi \\
& +\frac{1}{2}\left\{\left|\sigma_{S}\right|^{2} \partial_{s s} \varphi+\left|\sigma_{X}\right|^{2} \partial_{x x} \varphi+2\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{s x} \varphi\right\} .
\end{aligned}
$$

Corollary 3.8. Let $(y, z)$ be a solution of the PDE

$$
\begin{align*}
\mathfrak{a}(y)(t, x, s) & =-f(t, x, s, y(t, x, s), z(t, x, s))  \tag{3.16}\\
\left|\sigma_{S}\right|^{2} z(t, x, s) & =\left|\sigma_{S}\right|^{2} \partial_{s} y(t, x, s)+\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{x} y(t, x, s), \tag{3.17}
\end{align*}
$$

with terminal condition $y(T, .,)=.g(.,$.$) . Then the triplet (Y, Z, O)$, where

$$
Y_{t}=y\left(t, X_{t}, S_{t}\right), Z_{t}=z\left(t, X_{t}, S_{t}\right)
$$

and $\left(O_{t}\right)$ is given by (3.5) is a solution to the BSDE (3.1).

## 4 Explicit solution for Föllmer-Schweizer decomposition in the basis risk context

### 4.1 General considerations

We will discuss in this section the important Föllmer-Schweizer decomposition, denoted shortly F-S decomposition. It is a generalization of the well-known Galtchouk-KunitaWatanabe decomposition for martingales, to the more general case of semimartingales. Our task will consist in providing explicit expressions for the F-S decomposition in several situations. Let $S$ be a special semimartingale with canonical decomposition $S=M^{S}+V^{S}$. In the sequel we will convene that the space $L^{2}\left(M^{S}\right)$ consists of the predictable processes $\left(Z_{t}\right)_{t \in[0, T]}$ such that $\mathbb{E}\left[\int_{0}^{T}\left|Z_{s}\right|^{2} d\left\langle M^{S}\right\rangle_{s}\right]<\infty$ and $L^{2}\left(V^{S}\right)$ will denote the set of all predictable processes $\left(Z_{t}\right)_{t \in[0, T]}$ such that $\mathbb{E}\left[\int_{0}^{T}\left|Z_{s}\right|^{2} d\left\|V^{S}\right\|_{s}\right]<\infty$. The intersection of these two spaces is denoted

$$
\begin{equation*}
\Theta:=L^{2}\left(M^{S}\right) \cap L^{2}\left(V^{S}\right) . \tag{4.1}
\end{equation*}
$$

The Föllmer-Schweizer decomposition is defined as follows.
Definition 4.1. Let $h$ be a (possibly complex valued) square integrable $\mathcal{F}_{T}$-measurable random variable. We say that $h$ admits an $F-S$ decomposition with respect to $S$ if it can be written as

$$
\begin{equation*}
h=h_{0}+\int_{0}^{T} Z_{s} d S_{s}+O_{T}, \mathbb{P}-\text { a.s. } \tag{4.2}
\end{equation*}
$$

where $h_{0}$ is an $\mathcal{F}_{0}$-measurable r.v., $Z \in \Theta$ and $O=\left(O_{t}\right)_{t \in[0, T]}$ is a square integrable martingale, strongly orthogonal to $M^{S}$.

## Remark 4.2.

1) The notion of weak and strong orthogonality is discussed for instance in [30, Section 4.3] and [22, Section 1.4b]. Let $L$ and $N$ be two $\mathcal{F}_{t}$-local martingales, with null initial value. $L$ and $N$ are said to be strongly orthogonal if $L N$ is a local martingale. If $L$ and $N$ are locally square integrable, then they are strongly orthogonal if and only if $\langle L, N\rangle=0$. The definition of locally square integrable martingale is given for instance just before [30, Theorem 49 in Chapter 1].
2) The $F$-S decomposition makes also sense for complex valued square integrable random variable $h$. In that case the triplet $\left(h_{0}, Z, O\right)$ is generally complex.
3) If $h$ admits an $F$-S decomposition (4.2) then the complex conjugate $\bar{h}$ admits an $F-S$ decomposition given by

$$
\begin{equation*}
\bar{h}=\bar{h}_{0}+\int_{0}^{T} \bar{Z}_{s} d S_{s}+\bar{O}_{T}, \mathbb{P}-a . s \tag{4.3}
\end{equation*}
$$

The F-S decomposition has been extensively studied in the literature: sufficient conditions on the process $S$ were given so that every square integrable random variable has such a decomposition. A well-known condition ensuring the existence of such a decomposition is the so called structure condition (SC).

Definition 4.3. We say that a special semimartingale $S=V^{S}+M^{S}$ satisfies the structure condition $(S C)$ if there exists a predictable process $\alpha$ such that

1. $V_{t}^{S}=\int_{0}^{t} \alpha_{s} d\left\langle M^{S}\right\rangle_{s}$,
2. $\int_{0}^{T} \alpha_{s}^{2} d\left\langle M^{S}\right\rangle_{s}<\infty$ a.s.

The latter quantity plays a central role in the F-S decomposition. The associated process

$$
\begin{equation*}
K_{t}:=\int_{0}^{t} \alpha_{s}^{2} d\left\langle M^{S}\right\rangle_{s} \text { for } t \in[0, T] \tag{4.4}
\end{equation*}
$$

is called mean variance trade-off process.
Remark 4.4. [24] proved that, under (SC) and the additional condition that the process $K$ is uniformly bounded, the $F$-S decomposition of any real valued square integrable random variable exists and it is unique. More recent papers about the subject are [34], [7] and references therein.

This general decomposition refers to the process $S$ as underlying and it will be applied in the context of mean-variance hedging under basis risk, where $X$ is an observable price process of a non-traded asset.
As in previous sections, we consider a couple $(X, S)$ verifying the martingale problem (2.8), and we suppose Assumption 2.4 to be fulfilled. In the sequel we do not necessarily assume (SC) for $S$.

Definition 4.5. Let $h$ be a square integrable $\mathcal{F}_{T}$-measurable random variable. We say that $h$ admits a weak F-S decomposition with respect to $S$ if it can be written as

$$
\begin{equation*}
h=h_{0}+\int_{0}^{T} Z_{s} d S_{s}+O_{T}, \mathbb{P}-\text { a.s. } \tag{4.5}
\end{equation*}
$$

where $h_{0}$ is an $\mathcal{F}_{0}$-measurable r.v., $Z$ is a predictable process such that $\int_{0}^{T}\left|Z_{s}\right|^{2} d\left\langle M^{S}\right\rangle_{s}<\infty$ a.s., $\int_{0}^{T}\left|Z_{s}\right|^{2} d\left\|V^{S}\right\|_{s}<\infty$ a.s. and $O$ is a local martingale such that $\left\langle O, M^{S}\right\rangle=0$ with $O_{0}=0$.

Finding a weak F-S decomposition (4.5) $\left(h_{0}, Z, O\right)$ for some r.v. $h$ is equivalent to find a solution $(Y, Z, O)$ of the BSDE

$$
\begin{equation*}
Y_{t}=h-\int_{t}^{T} Z_{s} d S_{s}-\left(O_{T}-O_{t}\right) . \tag{4.6}
\end{equation*}
$$

The link is given by $Y_{0}=h_{0}$. The latter equation (4.6) can be seen as a special case of BSDE (3.1), where the driver $f$ is linear in $z$, of the form

$$
\begin{equation*}
f(t, x, s, y, z)=-\mathfrak{a}(i d)(t, x, s) z . \tag{4.7}
\end{equation*}
$$

This point of view was taken for instance by [33].
Remark 4.6. Let $(Y, Z, O)$ be a solution of (4.6) with $Z \in \Theta$, where $\Theta$ has been defined in (4.1) and $O$ is a square integrable martingale. Then $h$ admits an $F$ - $S$ decomposition (4.2) with $Y_{0}=h_{0}$.

We consider the case of the final value $h=g\left(X_{T}, S_{T}\right)$ for some continuous function $g$. Theorem 3.2 can be applied to obtain the result below.

Corollary 4.7. Let $y$ (resp. z): $[0, T] \times \mathcal{O} \rightarrow \mathbb{C}$. We suppose the following.

1) $y, \widetilde{y}:=y \times$ id belong to $\mathfrak{D}(\mathfrak{a})$.
2) $z$ verifies (3.15) of Remark 3.3.
3) $(y, z)$ solve the problem

$$
\begin{align*}
\mathfrak{a}(y)(t, x, s) & =\mathfrak{a}(i d)(t, x, s) z(t, x, s),  \tag{4.8}\\
\tilde{\mathfrak{a}}(y)(t, x, s) & =\widetilde{\mathfrak{a}}(i d)(t, x, s) z(t, x, s), \tag{4.9}
\end{align*}
$$

where the equalities hold in $\mathcal{L}$, with the terminal condition $y(T, .,)=.g(.,$.$) .$
Then the triplet $(Y, Z, O)$, where

$$
Y_{t}=y\left(t, X_{t}, S_{t}\right), Z_{t}=z\left(t, X_{t-}, S_{t-}\right), O_{t}=Y_{t}-Y_{0}-\int_{0}^{t} Z_{s} d S_{s}
$$

is a solution to the linear BSDE (4.6) linked to the weak $F$-S decomposition.

Remark 4.8. We remind that, setting $h_{0}=y\left(0, X_{0}, S_{0}\right)$, the triplet $\left(h_{0}, Z, O\right)$ is a candidate for a true F-S decomposition, see Definition 4.1. Sufficient conditions for this are the following.
a) $h=g\left(X_{T}, S_{T}\right) \in L^{2}(\Omega)$.
b) $\left(z\left(t, X_{t-}, S_{t-}\right)\right)_{t} \in \Theta$ i.e.

- $\mathbb{E}\left[\int_{0}^{T}\left|z\left(t, X_{t-}, S_{t-}\right)\right|^{2} \widetilde{\mathfrak{a}}(i d)\left(t, X_{t-}, S_{t-}\right) d A_{t}\right]<\infty$.
- $\mathbb{E}\left[\left(\int_{0}^{T}\left|z\left(t, X_{t-}, S_{t-}\right)\right|\left\|\widetilde{\mathfrak{a}}(i d)\left(t, X_{t-}, S_{t-}\right) d A\right\|_{t}\right)^{2}\right]<\infty$.
c) $\left(y\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(y)\left(u, X_{u-}, S_{u-}\right) d A_{u}\right)$ is an $\mathcal{F}_{t}$-square integrable martingale.

We remark that b) and c) imply by additivity that $O$ is a square integrable martingale. In fact

$$
\begin{equation*}
O_{t}=y\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(y)\left(u, X_{u-}, S_{u-}\right) d A_{u}-\int_{0}^{t} z\left(u, X_{u-}, S_{u-}\right) d M_{u}^{S}, \forall t \in[0, T] . \tag{4.10}
\end{equation*}
$$

### 4.2 Application: exponential of additive processes

We will investigate in this section a significant context where the equations in Corollary 4.7 can be solved, yielding the weak F-S decomposition and we can give sufficient conditions so that the true F-S decomposition is fulfilled. We focus on exponential of additive processes. Another example will be given in Section 4.3.

Let $(X, S)$ be a couple of exponential of semimartingale additive processes, as introduced in Section 2.5.

Proposition 4.9. Under Assumption 2.23, $S$ verifies the (SC) condition given in Definition 4.3 if and only if

$$
\begin{equation*}
\int_{0}^{T}\left(\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}\right)^{2} d \rho_{t}^{S}<\infty \text { a.s. } \tag{4.11}
\end{equation*}
$$

In this case, the mean variance trade-off process $K$ is deterministic and given by

$$
\begin{equation*}
K_{t}=\int_{0}^{t}\left(\frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}}\right)^{2} d \rho_{u}^{S}<\infty, \forall t \in[0, T] . \tag{4.12}
\end{equation*}
$$

Proof. It follows from Corollary 2.27 and item 5) of Proposition 2.24.
We look for the F-S decomposition of an $\mathcal{F}_{T}$-measurable random variable $h$ of the form $h:=g\left(X_{T}, S_{T}\right)$ for a function $g$ such that

$$
\begin{equation*}
g(x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \tag{4.13}
\end{equation*}
$$

where $\Pi$ is finite Borel complex measure.

In Section 2.5, Corollary 2.34 states that $(X, S)$ fulfills the martingale problem with respect to $\left(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, \rho^{S}\right)$ where the objects $\mathfrak{D}(\mathfrak{a}), \mathfrak{a}$ and $\rho^{S}$ were introduced respectively in (2.49), (2.50), (2.30). In order to determine the F-S decomposition (in its weak form given in (4.5)) we make use of Corollary 4.7. We look for a function $y$ (resp. $z$ ): $[0, T] \times \mathbb{R}^{2} \rightarrow \mathbb{C}$ such that Hypotheses 1), 2) and 3) are fulfilled. In agreement with definition of $\mathfrak{D}(\mathfrak{a})$ given in (2.49) we select $y$ of the form

$$
\begin{equation*}
y(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right), \tag{4.14}
\end{equation*}
$$

where $\Pi$ being the same finite complex measure as in (4.13) and $\lambda:[0, T] \times \mathbb{C}^{2} \rightarrow \mathbb{C}$. We will start by writing "necessary" conditions for a couple $(y, z)$, such that $y$ has the form (4.14), to be solutions of (4.8) and (4.9).

Suppose that the couple ( $y, z$ ) fulfills (4.8) and (4.9) of Corollary 4.7. We consider the expressions of $\mathfrak{a}(i d), \widetilde{\mathfrak{a}}(i d)$ given by (2.53), (2.54), and $\mathfrak{a}(y), \widetilde{\mathfrak{a}}(y)$ given by (2.50) and (2.55), for $f=y$. We replace them in the two above mentioned conditions (4.8) and (4.9) to obtain the following equations for $\lambda$ ( $d \rho_{t}^{S}$ a.e.).

$$
\begin{array}{rlr}
\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}}\left\{\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right\} & =s \frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}} z(t, x, s)  \tag{4.15}\\
\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) \lambda\left(t, z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}+1} \frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}} & = & s^{2} z(t, x, s) .
\end{array}
$$

The final condition $y(T, \cdot, \cdot)=g$ produces

$$
\begin{equation*}
\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \lambda\left(T, z_{1}, z_{2}\right)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \tag{4.16}
\end{equation*}
$$

Replacing $z$ from the second line of (4.15) in the first line, by identification of the inverse Fourier-Laplace transform, it follows that $\lambda$ verifies

$$
\begin{align*}
\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}} & =\lambda\left(t, z_{1}, z_{2}\right)\left\{\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}} \frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}}-\frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right\}  \tag{4.17}\\
\lambda\left(T, z_{1}, z_{2}\right) & =1, \tag{4.18}
\end{align*}
$$

for all $\left(z_{1}, z_{2}\right) \in \operatorname{supp} \Pi$. Without restriction of generality we can clearly set $\lambda\left(\cdot, z_{1}, z_{2}\right)=0$ for $\left(z_{1}, z_{2}\right)$ outside the support of $\Pi$. We observe that for fixed $z_{1}, z_{2}$, (4.17) constitutes an ordinary differential equation (in the Lebesgue-Stieltjes sense) in time $t$.

We solve now the linear differential equation (4.17). Provided that

$$
\begin{equation*}
u \mapsto \frac{d \rho_{u}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{u}^{S}} \frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}} \in L^{1}\left([0, T], d \rho_{s}\right), \tag{4.19}
\end{equation*}
$$

the (unique) solution of (4.17), is given by

$$
\begin{align*}
\lambda\left(t, z_{1}, z_{2}\right) & =\exp \left(\int_{t}^{T}\left[\frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}-\frac{d \rho_{u}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{u}^{S}} \frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}}\right] d \rho_{u}^{S}\right) \\
& =\exp \left(\int_{t}^{T} \kappa_{d u}\left(z_{1}, z_{2}\right)-\frac{d \rho_{u}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{u}^{S}} \kappa_{d u}(0,1)\right)  \tag{4.20}\\
& =\exp \left(\int_{t}^{T} \eta\left(z_{1}, z_{2}, d u\right)\right)
\end{align*}
$$

where

$$
\begin{equation*}
\eta\left(z_{1}, z_{2}, t\right):=\kappa_{t}\left(z_{1}, z_{2}\right)-\int_{0}^{t} \frac{d \rho_{u}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{u}^{S}} \kappa_{d u}(0,1) \tag{4.21}
\end{equation*}
$$

which is clearly absolutely continuous with respect to $d \rho^{S}$.
At this point, we have an explicit form of $\lambda$ defining the function $y$ intervening in the weak $\mathrm{F}-\mathrm{S}$ decomposition. In the sequel we will show that such a choice of $\lambda$ will constitute a sufficient condition so that $(y, z)$ where $y$ is defined by (4.14) and $z$ determined by the second line of (4.15), is a solution of the deterministic problem given by (4.8) and (4.9).

In order to check (4.19) and the validity of (4.15) and (4.16), we formulate the following assumption reinforcing Assumptions 2.23 and 2.31.

Assumption 4.10. Recall $I_{0}:=\operatorname{Re}(\operatorname{supp} \Pi)\left(\subset \mathbb{R}^{2}\right)$, where we convene that $\operatorname{Re}\left(z_{1}, z_{2}\right)=$ $\left(\operatorname{Re}\left(z_{1}\right), \operatorname{Re}\left(z_{2}\right)\right)$. We denote $I:=2 I_{0} \cup\{(0,1)\}$ and $\mathcal{D}$ the set

$$
\begin{equation*}
\mathcal{D}=\left\{z \in D, \int_{0}^{T}\left|\frac{d \kappa_{u}\left(z_{1}, z_{2}\right)}{d \rho_{u}^{S}}\right|^{2} d \rho_{u}^{S}<\infty\right\} \tag{4.22}
\end{equation*}
$$

We assume the validity of the properties below.

1) $\rho^{S}$ is strictly increasing.
2) $I_{0}$ is bounded.
3) $\forall z \in \operatorname{supp} \Pi, z, z+(0,1) \in \mathcal{D}$.
4) $\sup _{x \in I}\left\|\frac{d\left(\kappa_{t}(x)\right)}{d \rho_{t}^{S}}\right\|_{\infty}<\infty$.

## Remark 4.11.

1) Assumptions 2.23 and 2.31 are consequences of Assumption 4.10.
2) Taking into account Remark 2.33, we emphasize that, for the rest of this section, the statements would not change if we consider that the quantities integrated with respect to the measure $\Pi$ are null outside its support.
3) $I \subset \mathcal{D}$, in particular $(0,1) \in \mathcal{D}$ because of item 4) of Assumption 4.10.
4) By previous item and Proposition 4.9, $S$ verifies the (SC) condition and the mean variance trade-off process $K$ given by (4.12) is deterministic.
5) $I_{0} \subset D / 2$ (i.e. supp $\Pi \subset D / 2$ ). This follows again by item 4) of Assumption 4.10.

In the sequel we will introduce the following notation.

$$
\begin{equation*}
\gamma_{t}\left(z_{1}, z_{2}\right):=\frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}}, \forall\left(z_{1}, z_{2}\right) \in D / 2, t \in[0, T] . \tag{4.23}
\end{equation*}
$$

Similarly to [17, Lemma 3.28], we can show the upper bounds below.
Lemma 4.12. Under Assumption 4.10, we have the following.

1) Condition (4.19) is verified for $t \in[0, T],\left(z_{1}, z_{2}\right) \in \operatorname{supp} \Pi$.
2) There is a positive constant $c_{1}$, such that $d \rho_{s}^{S}$ a.e. $\sup _{\left(z_{1}, z_{2}\right) \in I_{0}+i \mathbb{R}^{2}} \frac{d \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, t\right)\right)}{d \rho_{t}^{S}} \leq c_{1}$.
3) There are positive constants $c_{2}, c_{3}$ such that, $d \rho_{s}$ a.e. the following holds.

$$
\text { For any }\left(z_{1}, z_{2}\right) \in I_{0}+i \mathbb{R}^{2},\left|\gamma_{t}\left(z_{1}, z_{2}\right)\right|^{2} \leq \frac{d \rho_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}} \leq c_{2}-c_{3} \frac{d \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, t\right)\right)}{d \rho_{t}^{S}}
$$

$$
\text { 4) } \sup _{\left(z_{1}, z_{2}\right) \in I_{0}+i \mathbb{R}^{2}}-\int_{0}^{T} 2 \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, d t\right)\right) \exp \left(\int_{t}^{T} 2 \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, d s\right)\right)\right)<\infty
$$

Proof. For illustration we prove item 1), the other points can be shown by similar techniques as in [17, Lemma 3.28].
Let $t \in[0, T],\left(z_{1}, z_{2}\right) \in \operatorname{supp} \Pi$. Condition (4.19) is valid since $(0,1) \in \mathcal{D}, z, z+(0,1) \in \mathcal{D}$ and

$$
\left(\int_{0}^{t}\left|\frac{d \rho_{u}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{u}^{S}} \frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}}\right| \rho_{d u}^{S}\right)^{2} \leq \int_{0}^{t}\left|\frac{d \rho_{u}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{u}^{S}}\right|^{2} \rho_{d u}^{S} \int_{0}^{t}\left|\frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}}\right|^{2} \rho_{d u}^{S}
$$

Now, we can state a proposition that gives indeed the weak F-S decomposition of a random variable $h=g\left(X_{T}, S_{T}\right)$.

Proposition 4.13. We suppose the validity of Assumption 4.10. Let $\lambda$ be defined as

$$
\begin{equation*}
\lambda\left(t, z_{1}, z_{2}\right)=\exp \left(\int_{t}^{T} \eta\left(z_{1}, z_{2}, d u\right)\right), \forall\left(z_{1}, z_{2}\right) \in D / 2 \tag{4.24}
\end{equation*}
$$

where $\eta$ has been defined at (4.21). Then $(Y, Z, O)$ is a solution of the BSDE (4.6), where

$$
\begin{aligned}
Y_{t} & =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) X_{t}^{z_{1}} S_{t}^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right) \\
Z_{t} & =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) X_{t-}^{z_{1}} S_{t-}^{z_{2}-1} \lambda\left(t, z_{1}, z_{2}\right) \gamma_{t}\left(z_{1}, z_{2}\right) \\
O_{t} & =Y_{t}-Y_{0}-\int_{0}^{t} Z_{s} d S_{s}
\end{aligned}
$$

recalling that $\gamma$ has been defined in (4.23).

Proof. The result will follow from Corollary 4.7 for which we need to check the assumptions. First we prove that the function $y$ defined by

$$
y(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right),
$$

where $\lambda$ is defined in (4.24), is indeed an element of $\mathfrak{D}(\mathfrak{a})$. Secondly, we prove that the associated $\widetilde{y}$ also belongs to $\mathfrak{D}(\mathfrak{a})$. Third, we check the condition (3.15) for $z$. Finally we need to check the validity of the system of equations (4.8) and (4.9).
Concerning $y$, the function $\lambda\left(\cdot, z_{1}, z_{2}\right)$ is well-defined for $\left(z_{1}, z_{2}\right) \in \operatorname{supp} \Pi$, thanks to point 1) of Lemma 4.12 and by definition we have $\lambda\left(d t, z_{1}, z_{2}\right) \ll \rho_{d t}^{S}, \quad \forall\left(z_{1}, z_{2}\right) \in D$, which is Condition (2.40).
In order to prove that $y \in \mathfrak{D}(\mathfrak{a})$, which was defined in (2.49), it remains to prove the two conditions below which constitute conditions (2.41) and (2.42) of Theorem 2.32.

$$
\begin{gather*}
\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2}<\infty, \quad \forall t \in[0, T] ;  \tag{4.25}\\
\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right|<\infty . \tag{4.26}
\end{gather*}
$$

Let $t \in[0, T],\left(z_{1}, z_{2}\right) \in D / 2$. By (4.24), we have

$$
\left|\lambda\left(t, z_{1}, z_{2}\right)\right|=\exp \left(\int_{t}^{T} \frac{d \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, u\right)\right)}{d \rho_{u}^{S}} \rho_{d u}^{S}\right),
$$

which implies, by item 2) of Lemma 4.12, that

$$
\begin{equation*}
\left|\lambda\left(t, z_{1}, z_{2}\right)\right| \leq \exp \left(c_{1} \rho_{T}^{S}\right), \tag{4.27}
\end{equation*}
$$

which gives in particular (4.25): in fact $\int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} \leq e^{2 c_{1} \rho_{T}^{S}}|\Pi|\left(\mathbb{C}^{2}\right)<\infty$.

Finally, to conclude that $y \in \mathfrak{D}(\mathfrak{a})$, we need to show (4.26). By construction, $\lambda$ verifies equation (4.17). Hence, by (4.17) and Cauchy-Schwarz we get

$$
\begin{align*}
\left(\int_{0}^{T} d \rho_{t}^{S} \left\lvert\, \frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right.\right. & \left.\left.+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}} \right\rvert\,\right)^{2}=\left(\int_{0}^{T} d \rho_{t}^{S}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|\left|\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}} \frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}}\right|\right)^{2} \\
& \leq \int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2}\left|\gamma_{t}\left(z_{1}, z_{2}\right)\right|^{2} d \rho_{t}^{S} \int_{0}^{T}\left|\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}\right|^{2} d \rho_{t}^{S}  \tag{4.28}\\
& \leq\left(I_{1}\left(z_{1}, z_{2}\right)+I_{2}\left(z_{1}, z_{2}\right)\right) \int_{0}^{T}\left|\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}\right|^{2} d \rho_{t}^{S}
\end{align*}
$$

with

$$
\begin{align*}
& I_{1}\left(z_{1}, z_{2}\right):=c_{2} \int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} d \rho_{t}^{S}, \\
& I_{2}\left(z_{1}, z_{2}\right):=-c_{3} \int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} \frac{d \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, t\right)\right)}{d \rho_{t}} d \rho_{t}^{S}, \tag{4.29}
\end{align*}
$$

where we have used item 3) of Lemma 4.12. Since $\lambda$ is uniformly bounded, see (4.27), we have

$$
\begin{equation*}
I_{1}\left(z_{1}, z_{2}\right) \leq c_{2} \rho_{T}^{S} \exp \left(2 c_{1} \rho_{T}^{S}\right) . \tag{4.30}
\end{equation*}
$$

On the other hand,

$$
\begin{align*}
I_{2}\left(z_{1}, z_{2}\right) & =-c_{3} \int_{0}^{T} \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, d t\right)\right) \exp \left(\int_{t}^{T} 2 \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, d s\right)\right)\right) \\
& \leq c_{3} \sup _{y \in I_{0}+i \mathbb{R}^{2}}-\int_{0}^{T} \operatorname{Re}\left(\eta\left(y_{1}, y_{2}, d t\right)\right) \exp \left(\int_{t}^{T} 2 \operatorname{Re}\left(\eta\left(y_{1}, y_{2}, d s\right)\right)\right), \tag{4.31}
\end{align*}
$$

which is finite by item 4 ) of Lemma 4.12. Integrating (4.28) with respect to $|\Pi|$, taking into account the two uniform bounds in $\left(z_{1}, z_{2}\right)$, i.e. (4.30) and (4.31), we can conclude to the validity of (4.26), so that $y \in \mathfrak{D}(\mathfrak{a})$.

We show similarly that $\widetilde{y}:=y \times i d \in \mathfrak{D}(\mathfrak{a})$. In fact, for $t \in[0, T]$ and $x, y>0$, we have

$$
\begin{aligned}
\widetilde{y}(t, x, s) & =\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}+1} \lambda\left(t, z_{1}, z_{2}\right) \\
& =\int_{\mathbb{C}^{2}} d \widetilde{\Pi}\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}} \widetilde{\lambda}\left(t, z_{1}, z_{2}\right),
\end{aligned}
$$

where $\widetilde{\lambda}\left(t, z_{1}, z_{2}\right)=\lambda\left(t, z_{1}, z_{2}-1\right)$ and $\widetilde{\Pi}$ is the Borel complex measure defined by

$$
\int_{\mathbb{C}^{2}} d \widetilde{\Pi}\left(z_{1}, z_{2}\right) \varphi\left(z_{1}, z_{2}\right)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) \varphi\left(z_{1}, z_{2}+1\right),
$$

for every bounded measurable function $\varphi$. Hence, $\operatorname{supp} \widetilde{\Pi}=\operatorname{supp} \Pi+(0,1)$. By Remark 4.11 1) and 5), we have $(0,1) \in D / 2$ and supp $\Pi \subset D / 2$. Then, by Remark 2.22 , supp $\widetilde{\Pi} \subset D$, so that Assumption 2.31 is verified for $\widetilde{\Pi}$. Moreover, by definition of $\widetilde{\Pi}$, the conditions (2.40) and (2.41) are fulfilled replacing $\Pi$ and $\lambda$ with $\widetilde{\Pi}$ and $\widetilde{\lambda}$. In order to conclude that $\widetilde{y} \in \mathfrak{D}(\mathfrak{a})$, we need to show

$$
\begin{equation*}
A:=\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\widetilde{\Pi}|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(t, z_{1}, z_{2}-1\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}-1\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right|<\infty, \tag{4.32}
\end{equation*}
$$

which corresponds to Condition (2.42) for $\Pi$ and $\lambda$ replaced by $\tilde{\Pi}$ and $\tilde{\lambda}$. Note that

$$
\begin{aligned}
A & =\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}+1\right)}{d \rho_{t}^{S}}\right| \\
& =\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right)\left(\frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}}+\frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}\right)\right| \\
& \leq A_{1}+A_{2}+A_{3},
\end{aligned}
$$

where

$$
\begin{aligned}
& A_{1}:=\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\frac{d \lambda\left(t, z_{1}, z_{2}\right)}{d \rho_{t}^{S}}+\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}}\right|, \\
& A_{2}:=\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right) \frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}\right|, \\
& A_{3}:=\int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right) \frac{d \rho_{t}\left(z_{1}, z_{2}, 0,1\right)}{d \rho_{t}^{S}}\right| .
\end{aligned}
$$

The first term $A_{1}$ is finite, since we already proved that $y \in \mathfrak{D}(\mathfrak{a})$ and so condition (4.26) is fulfilled. Moreover

$$
A_{2} \leq\left\|\frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}}\right\|_{\infty} \int_{0}^{T} d \rho_{t}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left|\lambda\left(t, z_{1}, z_{2}\right)\right|
$$

The right-hand side is finite, thanks to point 4) of Assumption 4.10 and the fact that $\lambda$ is uniformly bounded.

Finally, by Cauchy-Schwarz and item 3) of Lemma 4.12, taking into account Notation (4.23), by similar arguments as (4.28), we have

$$
\begin{aligned}
\left(A_{3}\right)^{2} & \leq|\Pi|\left(\mathbb{C}^{2}\right) \rho_{T}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \int_{0}^{T} d \rho_{t}^{S}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2}\left|\gamma_{t}\left(z_{1}, z_{2}\right)\right|^{2} \\
& \leq|\Pi|\left(\mathbb{C}^{2}\right) \rho_{T}^{S} \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right)\left(I_{1}\left(z_{1}, z_{2}\right)+I_{2}\left(z_{1}, z_{2}\right)\right)
\end{aligned}
$$

where $I_{1}\left(z_{1}, z_{2}\right)$ and $I_{2}\left(z_{1}, z_{2}\right)$ have been defined in (4.29). We have already shown in (4.30) and (4.31) that $I_{1}$ and $I_{2}$ are bounded on supp $\Pi$, hence $A_{3}<\infty$. In conclusion, it follows indeed that $\widetilde{y} \in \mathfrak{D}(\mathfrak{a})$ and Hypothesis 1) of Corollary 4.7 is verified.
We define $(t, x, s) \mapsto z(t, x, s)$ so that $s^{2} z(t, x, s)=\widetilde{\mathfrak{a}}(y)(t, x, s)$. This gives

$$
\begin{equation*}
z(t, x, s)=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) x^{z_{1}} s^{z_{2}-1} \lambda\left(t, z_{1}, z_{2}\right) \gamma_{t}\left(z_{1}, z_{2}\right), \forall t \in[0, T], x, y>0 \tag{4.33}
\end{equation*}
$$

Lemma 4.14 below shows that (3.15) is fulfilled and so Hypothesis 2) of Corollary 4.7 is verified.

We go on verifying Hypothesis 3) of Corollary 4.7, i.e. the validity of (4.15) and (4.16). Condition (4.16) is straightforward since $\lambda(T, \cdot, \cdot)=1$. The second equality in (4.15) takes place by definition of $z$. The first equality holds true integrating (4.17) thanks to (4.26). This proves 3) of Corollary 4.7.

Finally Corollary 4.7 implies that $(Y, Z, O)$, is a solution of the BSDE (4.6) provided we establish the following.

Lemma 4.14. Let $z$ be as in (4.33), where $\lambda, \gamma$ have been respectively defined in (4.24) and (4.23). We have

$$
\mathbb{E}\left[\int_{0}^{T}\left|z\left(r, X_{r-}, S_{r-}\right)\right|^{2} S_{r-}^{2} \rho_{d r}^{S}\right]<\infty
$$

In particular (3.15) is fulfilled.
Proof. First, let us show that

$$
\begin{equation*}
\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) \int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} \rho_{d t}\left(z_{1}, z_{2}\right)<\infty \tag{4.34}
\end{equation*}
$$

For this, we use points 3) and 4) of Lemma 4.12, (4.27) and (4.20) we get

$$
\begin{aligned}
\int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} \rho_{d t}\left(z_{1}, z_{2}\right)= & \int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} \frac{d \rho_{t}\left(z_{1}, z_{2}\right)}{d \rho_{t}^{S}} \rho_{d t}^{S} \\
\leq & \int_{0}^{T}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2}\left(c_{2}-c_{3} \frac{d \operatorname{Re}\left(\eta\left(y_{1}, y_{2}, t\right)\right)}{d \rho_{t}^{S}}\right) \rho_{d t}^{S} \\
\leq & c_{2} e^{2 c_{1} \rho_{T}^{S}} \rho_{T}^{S}-c_{3} \int_{0}^{T} \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, d t\right)\right) \exp \left(\int_{t}^{T} 2 \operatorname{Re}\left(\eta\left(z_{1}, z_{2}, d s\right)\right)\right) \\
\leq & c_{2} e^{2 c_{1} \rho_{T}^{S}} \rho_{T}^{S}+ \\
& c_{3} \sup _{\left(\xi_{1}, \xi_{2}\right) \in I_{0}+i \mathbb{R}^{2}}-\int_{0}^{T} \operatorname{Re}\left(\eta\left(\xi_{1}, \xi_{2}, d t\right)\right) \exp \left(\int_{t}^{T} 2 \operatorname{Re}\left(\eta\left(\xi_{1}, \xi_{2}, d s\right)\right)\right)
\end{aligned}
$$

Hence (4.34) is fulfilled.
Using Cauchy-Schwarz inequality, Fubini theorem and point 3) of Lemma 4.12, we have

$$
\begin{aligned}
& \mathbb{E}\left[\int_{0}^{T}\left|z\left(r, X_{r-}, S_{r-}\right)\right|^{2} S_{r-}^{2} d \rho_{r}^{S}\right]=\mathbb{E}\left[\int_{0}^{T}\left|\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) X_{t-}^{z_{1}} S_{t-}^{z_{2}} \lambda\left(t, z_{1}, z_{2}\right) \gamma_{t}\left(z_{1}, z_{2}\right)\right|^{2} \rho_{d s}^{S}\right] \\
\leq & |\Pi|\left(\mathbb{C}^{2}\right) \sup _{t \in[0, T],(a, b) \in I_{0}} \mathbb{E}\left[X_{t}^{2 a} S_{t}^{2 b}\right] \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \int_{0}^{t}\left|\lambda\left(t, z_{1}, z_{2}\right) \gamma_{t}\left(z_{1}, z_{2}\right)\right|^{2} \rho_{d s}^{S} \\
\leq & |\Pi|\left(\mathbb{C}^{2}\right) \sup _{t \in[0, T],(a, b) \in I_{0}} \mathbb{E}\left[X_{t}^{2 a} S_{t}^{2 b}\right] \int_{\mathbb{C}^{2}} d|\Pi|\left(z_{1}, z_{2}\right) \int_{0}^{t}\left|\lambda\left(t, z_{1}, z_{2}\right)\right|^{2} \rho_{d t}\left(z_{1}, z_{2}\right) .
\end{aligned}
$$

The right-hand side is finite, thanks to (4.34).

One can prove that the weak F-S decomposition in Proposition 4.13 is actually a strong F-S decomposition in the sense of Definition 4.1.

Theorem 4.15. Under Assumption 4.10, the random variable

$$
h=\int_{\mathbb{C}^{2}} d \Pi\left(z_{1}, z_{2}\right) X_{T}^{z_{1}} S_{T}^{z_{2}}
$$

admits an $F-S$ decomposition (4.2) where $h_{0}=Y_{0}$ and $(Y, Z, O)$ is given in Proposition 4.13.
Moreover, if $h$ is real-valued then the decomposition $(Y, Z, O)$ is real-valued and it is therefore the unique $F-S$ decomposition.

Remark 4.16. This statement is a generalization of the results of [17] (and [19]) to the case of hedging under basis risk. This yields a characterization of the hedging strategy in terms of Fourier-Laplace transform and the moment generating function.

Proof. Since $\Pi$ is a finite measure, then $h$ is square integrable. Indeed by Cauchy-Schwarz

$$
\begin{align*}
\mathbb{E}\left[h^{2}\right] & \leq|\Pi|\left(\mathbb{C}^{2}\right) \int_{\mathbb{C}^{2}} \mathbb{E}\left[\left|X_{T}\right|^{2 \operatorname{Re}\left(z_{1}\right)}\left|S_{T}\right|^{2 \operatorname{Re}\left(z_{2}\right)}\right] d \Pi\left(z_{1}, z_{2}\right)  \tag{4.35}\\
& \leq\left(|\Pi|\left(\mathbb{C}^{2}\right)\right)^{2} \sup _{(a, b) \in I} \mathbb{E}\left[\left|X_{T}\right|^{a}\left|S_{T}\right|^{b}\right]
\end{align*}
$$

where $I$ is a bounded subset of $\mathbb{R}^{2}$ defined in Assumption 4.10. By item 2) of Assumption 4.10 and item 3) of Proposition 2.24, previous quantity is finite.

By item 4) of Remark 4.11 and by Remark 4.4, the real-valued F-S decomposition of any real valued square integrable $\mathcal{F}_{T}$-measurable random variable is unique.

As a consequence, if $h$ is real-valued then its F-S decomposition is also real-valued. In fact, if $\left(Y_{0}, Z, O\right)$ is an F-S decomposition of $h$, then $\left(\bar{Y}_{0}, \bar{Z}, \bar{O}\right)$ is also an F-S of $\bar{h}$ by item 3) of Remark 4.2. Thus, by subtraction, $\left(\operatorname{Im}\left(Y_{0}\right), \operatorname{Im}(Z), \operatorname{Im}(O)\right)$ is an F-S decomposition with real-valued triplet of the real-valued r.v. $\operatorname{Im}(h)=0$. By uniqueness $\operatorname{Im}\left(Y_{0}\right), \operatorname{Im}(Z)$ and $\operatorname{Im}(O)$ are null and the decomposition $\left(Y_{0}, Z, O\right)$ is real valued.

Now, let $(Y, Z, O)$ defined in Proposition 4.13. It remains to prove that $\left(Y_{0}, Z, O\right)$ is a strong (possibly complex) F-S decomposition in the sense of Definition 4.1. For this we need to show items a), b), c) of Remark 4.8. Item a) has been the object of (4.35).

We show below item b) i.e. $\mathbb{E}\left[\int_{0}^{T}\left|Z_{s}\right|^{2} d\left\langle M^{S}\right\rangle_{s}\right]<\infty$ and $\mathbb{E}\left[\left(\int_{0}^{T}\left|Z_{s}\right| d\left\|V^{S}\right\|_{s}\right)^{2}\right]<\infty$. The first inequality is stated in Lemma 4.14. In order to prove the second one, we remind that, by Corollary 2.27,

$$
d V_{t}^{S}=S_{t-} \kappa_{d t}(0,1)=S_{t-} \frac{d \kappa_{t}(0,1)}{d \rho_{t}^{S}} \rho_{d t}^{S}
$$

Consequently

$$
\begin{aligned}
\mathbb{E}\left[\left(\int_{0}^{T}\left|Z_{s}\right| d\left\|V^{S}\right\|_{s}\right)^{2}\right] & =\mathbb{E}\left[\left(\int_{0}^{T}\left|Z_{u}\right|\left|\frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}}\right| S_{u-} \rho_{d u}^{S}\right)^{2}\right] \\
& \leq \int_{0}^{T}\left|\frac{d \kappa_{u}(0,1)}{d \rho_{u}^{S}}\right|^{2} \rho_{d u}^{S} \mathbb{E}\left[\int_{0}^{T}\left|Z_{u}\right|^{2} S_{u-}^{2} \rho_{d u}^{S}\right]
\end{aligned}
$$

which is finite since, by item 3 ) of Remark 4.11 which says that $(0,1) \in \mathcal{D}$, taking into account Lemma 4.14.

To end this proof, we need to show item c) of Remark 4.8. For this we use Proposition 2.35 for which we need to check conditions a) and b). By item 5) of Remark 4.11 we have $I_{0} \subset D / 2$ which constitutes item a). Item b) is verified by condition (4.34) is verified. Hence Proposition 2.35 implies that

$$
t \mapsto y\left(t, X_{t}, S_{t}\right)-\int_{0}^{t} \mathfrak{a}(y)\left(u, X_{u-}, S_{u-}\right) \rho_{d t}^{S}
$$

is a square integrable martingale.

### 4.3 Diffusion processes

We set $\mathcal{O}=\mathbb{R} \times E$, where $E=\mathbb{R}$ or $] 0, \infty[$. In this Section we apply Corollary 4.7 to the diffusion processes $(X, S)$ modeled in Section 2.3 whose dynamics is given by (2.28). We are interested in the F-S decomposition of $h=g\left(X_{T}, S_{T}\right)$. We recall the assumption in that context.

## Assumption 4.17.

- $b_{X}, b_{S}, \sigma_{X}$ and $\sigma_{S}$ are continuous and globally Lipschitz.
- $g: \mathcal{O} \rightarrow \mathbb{R}$ is continuous.

We remind that $(X, S)$ solve the strong martingale problem related to $(\mathfrak{D}(\mathfrak{a}), \mathfrak{a}, A)$ where $A_{t}=t, \mathfrak{D}(\mathfrak{a})=\mathcal{C}^{1,2}\left(\left[0, T[\times \mathcal{O}) \cap \mathcal{C}^{1}([0, T] \times \mathcal{O})\right.\right.$. For a function $y \in \mathfrak{D}(\mathfrak{a})$, obviously $\tilde{y} \in \mathfrak{D}(\mathfrak{a})$ and the operators $\mathfrak{a}$ and $\widetilde{\mathfrak{a}}$ are given by

$$
\begin{aligned}
\mathfrak{a}(y) & =\partial_{t} y+b_{S} \partial_{s} y+b_{X} \partial_{x} y \\
& +\frac{1}{2}\left\{\left|\sigma_{S}\right|^{2} \partial_{s s} y+\left|\sigma_{X}\right|^{2} \partial_{x x} y+2\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{s x} y\right\}, \\
\tilde{\mathfrak{a}}(y) & =\left|\sigma_{S}\right|^{2} \partial_{s} y+\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{x} y .
\end{aligned}
$$

Conditions 3 ) of that Corollary 4.7 translates into

$$
\begin{align*}
b_{S} z & =\partial_{t} y+b_{S} \partial_{s} y+b_{X} \partial_{x} y+\frac{1}{2}\left\{\left|\sigma_{S}\right|^{2} \partial_{s s} y+\left|\sigma_{X}\right|^{2} \partial_{x x} y+2\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{s x} y\right\} \\
y(T, ., .) & =g(., .)  \tag{4.36}\\
\left|\sigma_{S}\right|^{2} z & =\left|\sigma_{S}\right|^{2} \partial_{s} y+\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{x} y .
\end{align*}
$$

If, moreover, $\frac{1}{\left|\sigma_{S}\right|}$ is locally bounded, then we have the following:

$$
\left\{\begin{array}{r}
\partial_{t} y+B \partial_{x} y+\frac{1}{2}\left(\left|\sigma_{S}\right|^{2} \partial_{s s} y+\left|\sigma_{X}\right|^{2} \partial_{x x} y+2\left\langle\sigma_{S}, \sigma_{X}\right\rangle \partial_{s x} y\right)=0  \tag{4.37}\\
y(T, ., .)=g(., .)
\end{array}\right.
$$

and

$$
\begin{equation*}
z=\partial_{s} y+\frac{\left\langle\sigma_{S}, \sigma_{X}\right\rangle}{\left|\sigma_{S}\right|^{2}} \partial_{x} y \tag{4.38}
\end{equation*}
$$

where

$$
\begin{equation*}
B=b_{X}-b_{S} \frac{\left\langle\sigma_{S}, \sigma_{X}\right\rangle}{\left|\sigma_{S}\right|^{2}} \tag{4.39}
\end{equation*}
$$

$z$ is then locally bounded since $\sigma_{S}, \sigma_{X}$ and $\frac{1}{\left|\sigma_{S}\right|}$ are locally bounded and because $y \in \mathfrak{D}(\mathfrak{a})$.
Proposition 4.18. We suppose the validity of Assumption 4.17 and that $\left|\sigma_{S}\right|$ is always strictly positive.
If $(y, z)$ is a solution of the system (4.37) and (4.38), such that $y \in \mathfrak{D}(\mathfrak{a})$, then $(Y, Z, O)$ is a solution of the BSDE (4.6), where

$$
\begin{aligned}
Y_{t} & =y\left(t, X_{t}, S_{t}\right), \\
Z_{t} & =z\left(t, X_{t}, S_{t}\right), \\
O_{t} & =Y_{t}-Y_{0}-\int_{0}^{t} Z_{s} d S_{s}
\end{aligned}
$$

Proof. It follows from Corollary 4.7 for which we need to check the conditions 1), 2) and 3). Indeed, since $y, \tilde{y} \in \mathfrak{D}(\mathfrak{a})$, Condition 1 ) holds; since $z$ is locally bounded, by item 2 . of Remark 3.3, Condition 2) is fulfilled. Condition 3) has been the object of the considerations above the statement of the Proposition.

The result above yields the weak F-S decomposition for $h$. In order to show that $\left(Y_{0}, Z, O\right)$ constitutes a true F-S decomposition, we need to make use of Remark 4.8. First we introduce the following assumption.

Assumption 4.19. Suppose that the process $(X, S)$ takes values in $\mathcal{O}$ and the following.
i) $g \in C^{1}$ such that $g, \partial_{x} g$ and $\partial_{s} g$ have polynomial growth.
ii) $B$ is globally Lipschitz.
iii) $\partial_{x} B, \partial_{s} B, \partial_{x} \sigma_{X}, \partial_{s} \sigma_{X}, \partial_{x} \sigma_{S}$ and $\partial_{s} \sigma_{S}$ exist, are continuous and have polynomial growth.
iv) $\sigma_{S}$ never vanishes.

We formulate the following.
Theorem 4.20. Suppose that Assumptions 4.17 and 4.19 are fulfilled, and suppose the existence of a function $y:[0, T] \times \mathcal{O} \rightarrow \mathbb{R}$ such that
$y \in C^{0}([0, T] \times \mathcal{O}) \cap C^{1,2}([0, T[\times \mathcal{O})$ verifies the $P D E(4.37)$ and has polynomial growth.

Then the $F$ - $S$ decomposition (4.2) of $h=g\left(X_{T}, S_{T}\right)$ is provided by $\left(h_{0}, Z, O\right)$ where, $h_{0}=Y_{0}$ and

$$
Y_{t}=y\left(t, X_{t}, S_{t}\right), Z_{t}=z\left(t, X_{t}, S_{t}\right), O_{t}=Y_{t}-Y_{0}-\int_{0}^{t} Z_{s} d S_{s}
$$

and $z:[0, T] \times \mathcal{O} \rightarrow \mathbb{R}$ is given by (4.38).
Proof. Let $y:[0, T] \times \mathcal{O} \rightarrow \mathbb{R}$ verifying (4.40) and $z$ defined by (4.38). In order to show that the triplet given in Proposition 4.18 yields a true F-S decomposition, we need to show items a), b), c) of Remark 4.8.

First note that the random variable $g\left(X_{T}, S_{T}\right)$ is square integrable, because $g$ has polynomial growth and $X$ and $S$ admit all moments, see Remark 2.20. So a) is verified.

In view of verifying item b) of Remark 4.8 we remind that

$$
\mathfrak{a}(i d)=b_{S}, \tilde{\mathfrak{a}}(i d)=\left|\sigma_{S}\right|^{2}, A_{t} \equiv t \text { and } z=\partial_{s} y+\frac{\left\langle\sigma_{S}, \sigma_{X}\right\rangle}{\left|\sigma_{S}\right|^{2}} \partial_{x} y
$$

Indeed, since $y$ has polynomial growth, it is forced to be unique since [23, Theorem 7.6 , chapter 5] implies that

$$
\begin{equation*}
y(t, x, s)=\mathbb{E}\left[g\left(X_{T}^{t, x, s}, S_{T}^{t, x, s}\right)\right] \tag{4.41}
\end{equation*}
$$

where $\left(\widetilde{X}=X^{t, x, s}, \widetilde{S}=S^{t, x, s}\right)$ is a solution of

$$
d\binom{\widetilde{X}_{r}}{\widetilde{S}_{r}}=\Sigma\left(r, \widetilde{X}_{r}, \widetilde{S}_{r}\right) d \widetilde{W}_{r}+\binom{B\left(r, \widetilde{X}_{r}, \widetilde{S}_{r}\right)}{0} d r
$$

with $\widetilde{X}_{t}=x, \widetilde{S}_{t}=s$, where $\widetilde{W}=\left(\widetilde{W}^{1}, \widetilde{W}^{2}\right)$ is a standard two-dimensional Brownian motion, and

$$
\Sigma=\left(\begin{array}{ll}
\sigma_{X, 1} & \sigma_{X, 2} \\
\sigma_{S, 1} & \sigma_{S, 2}
\end{array}\right)
$$

We remind that $B$ has been defined in (4.39).
By (4.41), a straightforward adaptation of [15, Theorem 5.5] yields that the partial derivatives $\partial_{x} y$ and $\partial_{s} y$ exist and are continuous on $[0, T] \times \mathcal{O}$ and they have polynomial growth.

Using (4.38), we have

$$
z b_{S}=b_{S} \partial_{s} y+b_{X} \partial_{x} y-B \partial_{x} y
$$

Now, since $\partial_{x} y$ and $\partial_{s} y$ have polynomial growth, and by assumption $b_{S}, b_{X}$ and $B$ have linear growth, we get that $z b_{S}$ has polynomial growth. This gives, by Remark 2.20,

$$
\mathbb{E}\left[\left(\int_{0}^{T}\left|z b_{S}\right|\left(t, X_{t}, S_{t}\right) d t\right)^{2}\right]<\infty
$$

On the other hand, using (4.38) and Cauchy-Schwarz, we have

$$
\begin{aligned}
\left|z \sigma_{S}\right| & =\| \sigma_{S}\left|\partial_{s} y+\frac{\left\langle\sigma_{X}, \sigma_{S}\right\rangle}{\left|\sigma_{S}\right|} \partial_{x} y\right| \\
& \leq\left|\sigma_{S} \| \partial_{s} y\right|+\left|\sigma_{X}\right|\left|\partial_{x} y\right|
\end{aligned}
$$

Since $\sigma_{X}, \sigma_{S}$ have linear growth and $\partial_{x} y$ and $\partial_{s} y$ have polynomial growth, we get that $z \sigma_{S}$ has polynomial growth, which implies, by Remark 2.20, that

$$
\mathbb{E}\left[\int_{0}^{T}\left|z \sigma_{S}\right|^{2}\left(t, X_{t}, S_{t}\right) d t\right]<\infty
$$

Consequently, item b) of Remark 4.8 is fulfilled.
In order to show the last item c), taking into account Remark 2.21, we need to prove that

$$
\begin{aligned}
u \mapsto M_{u}^{Y} & =\int_{0}^{u} \partial_{x} y\left(r, X_{r}, S_{r}\right)\left(\sigma_{X, 1}\left(r, X_{r}, S_{r}\right) d W_{r}^{1}+\sigma_{X, 2}\left(r, X_{r}, S_{r}\right) d W_{r}^{2}\right) \\
& +\int_{0}^{u} \partial_{s} y\left(r, X_{r}, S_{r}\right)\left(\sigma_{S, 1}\left(r, X_{r}, S_{r}\right) d W_{r}^{1}+\sigma_{S, 2}\left(r, X_{r}, S_{r}\right) d W_{r}^{2}\right)
\end{aligned}
$$

is a square integrable martingale. This is due to the fact that $\partial_{x} y$ and $\partial_{s} y$ have polynomial growth, and that $\sigma_{X}$ and $\sigma_{S}$ have linear growth, and Remark 2.20 , which implies that

$$
\mathbb{E}\left[\int_{0}^{T}\left\{\left(\partial_{x} y\left(r, X_{r}, S_{r}\right)\right)^{2}\left|\sigma_{X}\left(r, X_{r}, S_{r}\right)\right|^{2}+\left(\partial_{s} y\left(r, X_{r}, S_{r}\right)\right)^{2}\left|\sigma_{S}\left(r, X_{r}, S_{r}\right)\right|^{2}\right\} d u\right]<\infty
$$

This concludes the proof of Theorem 4.20.
Below we show that, under Assumptions 4.17 and 4.19 , Condition (4.40) is not really restrictive.

Proposition 4.21. We assume the validity of Assumptions 4.17 and 4.19.
Moreover we suppose the validity of one of the three items below.

1) We set $\mathcal{O}=\mathbb{R}^{2}$. Suppose that the second (partial, with respect to $(x, s)$ ) derivatives of $B, \sigma_{X}, \sigma_{S}$ and $g$ exist, are continuous and have polynomial growth.
2) We set $\mathcal{O}=\mathbb{R}^{2}$. We suppose $B, \sigma_{X}, \sigma_{S}$ to be bounded and there exist $\lambda_{1}, \lambda_{2}>0$ such that

$$
\lambda_{1}|\xi|^{2} \leq\left(\xi_{1}, \xi_{2}\right) C(t, x, s)\left(\xi_{1}, \xi_{2}\right)^{T} \leq \lambda_{2}|\xi|^{2}, \forall \xi=\left(\xi_{1}, \xi_{2}\right) \in \mathcal{O},
$$

where $C(t, x, s)=\left(\left.\begin{array}{c}\left|\sigma_{X}\right|^{2}(t, x, s) \\ \left\langle\sigma_{X}, \sigma_{S}\right\rangle(t, x, s)\end{array} \right\rvert\, \begin{array}{c}\left\langle\sigma_{X}, \sigma_{S}\right\rangle(t, x, s) \\ \left|\sigma_{S}\right|^{2}(t, x, s)\end{array}\right)$.
3) (Black-Scholes case.) We suppose $\mathcal{O}=] 0,+\infty\left[{ }^{2}\right.$.

$$
\begin{aligned}
b_{S}(t, x, s) & =s \hat{b}_{S}, \\
b_{X}(t, x, s) & =x \hat{b}_{X} \\
\sigma_{S}(t, x, s) & =\left(s \hat{\sigma}_{S, 1}, s \hat{\sigma}_{S, 2}\right) \\
\sigma_{X}(t, x, s) & =\left(x \hat{\sigma}_{X, 1}, x \hat{\sigma}_{X, 2}\right),
\end{aligned}
$$

where $\hat{b}_{S}, \hat{b}_{X}, \hat{\sigma}_{S, 1}, \hat{\sigma}_{S, 2}, \hat{\sigma}_{X, 1}$ and $\hat{\sigma}_{X, 2}$ are constants, such that $\left\langle\hat{\sigma}_{X}, \hat{\sigma}_{S}\right\rangle<\left|\hat{\sigma}_{X}\right|\left|\hat{\sigma}_{S}\right|$.
We have the following.
i) There is a (unique) strict solution $y$ of (4.37) in the class $C^{1,2}\left(\left[0, T[\times \mathcal{O}) \cap C^{0}([0, T] \times \mathcal{O})\right.\right.$ with polynomial growth.
ii) The F-S decomposition (4.2) of $h=g\left(X_{T}, S_{T}\right)$ is provided by $\left(h_{0}, Z, O\right)$ where $(Y, Z, O)$ fulfill

$$
Y_{t}=y\left(t, X_{t}, S_{t}\right), Z_{t}=z\left(t, X_{t}, S_{t}\right) \text { and } O_{t}=Y_{t}-Y_{0}-\int_{0}^{t} Z_{s} d S_{s},
$$

where $z$ is given by (4.38).
Remark 4.22. We will show below that under the hypotheses above, conclusion i) holds, i.e. there is a function $y$ fulfiling (4.40). We observe that, by the proof of Theorem 4.20, if such a $y$ exists then it admits the probabilistic representation (4.41) and so it is necessarily the unique $C^{1,2}\left(\left[0, T[\times \mathcal{O}) \cap C^{0}([0, T] \times \mathcal{O})\right.\right.$, with polynomial growth, solution of (4.37).

Proof. We proceed to discussing the existence of $y$ mentioned in Remark 4.22. So we distinguish now the mentioned three cases.

Suppose first item 1). The function $y$ defined by (4.41) is a continuous function by the fact that the flow ( $\widetilde{X}, \widetilde{S}$ ) is continuous in all variables and Remark 2.20, taking into account Lebesgue dominated convergence theorem. [15, Theorem 6.1], states that $y$ belongs to $C^{1,2}([0, T] \times \mathcal{O})$, and it verifies the PDE (4.37). [15, Theorem 5.5$]$ says in particular that $y$ has polynomial growth. In that case conclusion i) is established.

Under the assumption described in item 2), the conclusion i) can be obtained by simply adapting the proof of [14, Theorem 12, p.25]. Indeed, according to [14, Theorem 8, p.19]
there is a fundamental solution $\Gamma:\left\{\left(t_{1}, t_{2}\right), 0 \leq t_{1}<t_{2} \leq T\right\} \times \mathbb{R}^{2} \times \mathbb{R}^{2} \rightarrow \mathbb{R}$ such that

$$
\begin{equation*}
\Gamma\left(t_{1}, t_{2} ; \gamma, \xi\right) \leq \frac{1}{a_{1}\left(t_{2}-t_{1}\right)} \exp \left(-\frac{-|\gamma-\xi|^{2}}{a_{1}\left(t_{2}-t_{1}\right)}\right) \tag{4.42}
\end{equation*}
$$

where $a_{1}$ is a positive constant.
Now, by [14, Theorem 12, p.25], the function $y$ defined by

$$
\begin{equation*}
y(t, x, s)=\int_{\mathbb{R}^{2}} \Gamma\left(t, T ;(x, s),\left(\xi_{1}, \xi_{2}\right)\right) g\left(\xi_{1}, \xi_{2}\right) d \xi_{1} d \xi_{2} \tag{4.43}
\end{equation*}
$$

is a strict solution of (4.37), in particular it belongs to $\mathcal{C}^{1,2}\left(\left[0, T\left[\times \mathbb{R}^{2}\right) \cap \mathcal{C}^{0}\left([0, T] \times \mathbb{R}^{2}\right)\right.\right.$.
Since g has polynomial growth then there exist $a_{2}>0, p>1$ such that, $\forall x, s \in \mathbb{R}$,

$$
\begin{equation*}
|g(x, s)| \leq a_{2}\left(1+|x|^{p}+|s|^{p}\right) \tag{4.44}
\end{equation*}
$$

Thus, by (4.43), (4.42) and (4.44), for $x, s \in \mathbb{R}$ and $0 \leq t \leq T$, we have

$$
|y(t, x, s)| \leq \frac{a_{2}}{a_{1}(T-t)} \int_{\mathbb{R}^{2}}\left(1+\left|\xi_{1}\right|^{p}+\left|\xi_{2}\right|^{p}\right) \exp \left(-\frac{\left|x-\xi_{1}\right|^{2}+\left|s-\xi_{2}\right|^{2}}{a_{1}(T-t)}\right) d \xi_{1} d \xi_{2}
$$

So there is a constant $C_{1}(p, T)>0$ such that

$$
\begin{equation*}
|y(t, x, s)| \leq C_{1}(p, T)\left(1+\mathbb{E}\left[\left|x+G_{1}\right|^{p}+\left|x+G_{2}\right|^{p}\right]\right) \tag{4.45}
\end{equation*}
$$

where $G=\left(G_{1}, G_{2}\right)$ is a two dimensional centered Gaussian vector with covariance matrix equal to $\frac{a_{1}(T-t)}{2}$ times the identity matrix. Since $p>1$, then there is a constant $C_{2}(p, T)$ such that

$$
\begin{aligned}
|y(t, x, s)| & \leq C_{2}(p, T)\left(1+|x|^{p}+|s|^{p}+\mathbb{E}\left[\left|G_{1}\right|^{p}+\left|G_{2}\right|^{p}\right]\right) \\
& \leq C_{3}(p, T)\left(1+|x|^{p}+|s|^{p}\right)
\end{aligned}
$$

where $C_{3}(p, T)$ is another positive constant. In conclusion the solution $y$ given by (4.43) has polynomial growth.

We discuss now the Black-Scholes case 3) showing that, also in that case, there is $y$ such that (4.40) is fulfilled. First note that the uniform ellipticity condition in 2) is not fulfilled for this dynamics, so we consider a logarithmic change of variable. For a function $y \in \mathfrak{D}(\mathfrak{a})$, we introduce the function $\hat{y}:[0, T] \times \mathbb{R}^{2} \rightarrow \mathbb{R}$ defined by

$$
\hat{y}(t, x, s)=y(t, \log (x), \log (s)), \forall t \in[0, T], x, s>0
$$

By inspection we can show that $y$ is a solution of (4.37) if and only if $\hat{y}$ fulfills

$$
\begin{align*}
0= & \partial_{t} \hat{y}+\left(\hat{b}_{X}-\hat{b}_{S} \frac{\left\langle\hat{\sigma}_{S}, \hat{\sigma}_{X}\right\rangle}{\left|\hat{\sigma}_{S}\right|^{2}}-\frac{1}{2}\left|\hat{\sigma}_{X}\right|^{2}\right) \partial_{x} \hat{y}-\frac{1}{2}\left|\hat{\sigma}_{S}\right|^{2} \partial_{s} \hat{y}+ \\
& +\frac{1}{2}\left(\left|\hat{\sigma}_{S}\right|^{2} \partial_{s s} \hat{y}+\left|\hat{\sigma}_{X}\right|^{2} \partial_{x x} \hat{y}+2\left\langle\hat{\sigma}_{S}, \hat{\sigma}_{X}\right\rangle \partial_{s x} \hat{y}\right)  \tag{4.46}\\
\hat{y}(T, ., .)= & \hat{g}(., .)
\end{align*}
$$

where $\hat{g}(x, s)=g\left(e^{x}, e^{s}\right), \forall x, s \in \mathbb{R}$. Note that the PDE problem (4.46) has constant coefficients and it verifies the uniform ellipticity condition in 2 ).

Moreover, since $g$ has polynomial growth, then there exist $c>0, p>1$ such that $g(x, s) \leq$ $c\left(1+x^{p}+s^{p}\right), \forall x, s>0$ again. Hence $\hat{g}(x, s) \leq c\left(1+e^{p x}+e^{p s}\right), \forall x, s \in \mathbb{R}$.

Again, by simple adaptation of the proof of [14, Theorem 12, p.25], we observe that equation (4.46) admits a solution $\hat{y}$ in $C^{1,2}\left(\left[0, T\left[\times \mathbb{R}^{2}\right) \cap C^{0}\left([0, T] \times \mathbb{R}^{2}\right)\right.\right.$, such that

$$
\hat{y}(t, x, s) \leq K\left(1+e^{p x}+e^{p s}\right), \forall x, s \in \mathbb{R}
$$

where $K>0$. This yields that $y$ has polynomial growth, since $y(t, x, s)=\hat{y}(t, \log (x), \log (s)), \forall t \in$ $[0, T], x, s>0$, so

$$
y(t, x, s) \leq K\left(1+x^{p}+s^{p}\right), \forall t \in[0, T], x, s>0
$$

This concludes the proof of conclusion i).
Conclusion ii) is now a direct consequence of Theorem 4.20 together with condition i).
Remark 4.23. The last item of Proposition 4.21 permits to recover the results already found in [20], by replacing

$$
\begin{aligned}
\hat{b}_{S} & =\left(\mu_{S}-r\right) \\
\hat{b}_{X} & =\left(\mu_{U}-r\right) \\
\hat{\sigma}_{S} & =\left(\sigma_{S}, 0\right) \\
\hat{\sigma}_{X} & =\left(\rho \sigma_{U}, \sqrt{1-\rho^{2}} \sigma_{U}\right)
\end{aligned}
$$

where $\mu_{S}, \mu_{U}, r, \sigma_{S}$ and $\sigma_{U}$ are constants.

## Appendix A Proof of Proposition 2.8

Proof. Let $f \in E$ and set $\tilde{f}(x)=\frac{f(x)}{1+x^{2}}, \forall x \in \mathbb{R}$. Condition (2.16) implies, by mean value theorem, that there exists a constant $c(t)$ such that

$$
\mathbb{E}\left[\left|X_{t}^{0, x}-X_{t}^{0, y}\right|^{2}\right] \leq c(t)|x-y|^{2}, \forall x, y \in \mathbb{R}
$$

Then, by the Garsia-Rodemich-Rumsey criterion, see for instance [2, Section 3], there exists a r.v. $\Gamma_{t}$ such that $\mathbb{E}\left[\Gamma_{t}^{2}\right]<\infty$ and $\forall x, y \in \mathbb{R}$

$$
\begin{equation*}
\left|X_{t}^{0, x}-X_{t}^{0, y}\right| \leq \Gamma_{t}|x-y|^{\alpha}, \text { for } 0<\alpha<\frac{1}{2} \tag{A.1}
\end{equation*}
$$

possibly up to a modified version of the flow.
This implies in particular that for $x \in \mathbb{R}$

$$
\begin{aligned}
\frac{\left|X_{t}^{0, x}\right|^{2}}{1+x^{2}} & \leq \frac{2}{1+x^{2}}\left(\left|X_{t}^{0,0}\right|^{2}+\left|X_{t}^{0, x}-X_{t}^{0,0}\right|^{2}\right) \\
& \leq \frac{2}{1+x^{2}}\left(\left|X_{t}^{0,0}\right|^{2}+\left|\Gamma_{t}\right|^{2}|x|^{2 \alpha}\right) \\
& \leq 2\left(\left|X_{t}^{0,0}\right|^{2}+\left|\Gamma_{t}\right|^{2}\right)
\end{aligned}
$$

Hence

$$
\begin{equation*}
\sup _{x \in \mathbb{R}} \mathbb{E}\left[\frac{\left|X_{t}^{0, x}\right|^{2}}{1+x^{2}}\right]<\infty . \tag{A.2}
\end{equation*}
$$

Consequently, for $x \in \mathbb{R}$, we have

$$
\begin{aligned}
\frac{\left|P_{t} f(x)\right|}{1+x^{2}} & =\frac{\left|\mathbb{E}\left[f\left(X_{t}^{0, x}\right)\right]\right|}{1+x^{2}} \\
& \leq\|f\|_{E} \frac{1+\mathbb{E}\left[\left|X_{t}^{0, x}\right|^{2}\right]}{1+x^{2}} \\
& \leq\|f\|_{E} \sup _{\xi \in \mathbb{R}} \frac{1+\mathbb{E}\left[\left|X_{t}^{0, \xi}\right|^{2}\right]}{1+\xi^{2}} .
\end{aligned}
$$

The right-hand side is finite, thanks to (A.2), so that

$$
\begin{equation*}
\left\|P_{t} f\right\|_{E} \leq\|f\|_{E} \sup _{\xi \in \mathbb{R}} \frac{1+\mathbb{E}\left[\left|X_{t}^{0, \xi}\right|^{2}\right]}{1+\xi^{2}} \tag{A.3}
\end{equation*}
$$

After we will have shown that $\widetilde{P_{t} f}$ is also uniformly continuous, (A.3) will also imply that $P_{t} f \in E$ and that $P_{t}$ is a bounded linear operator.

Therefore it remains to show that $\widetilde{P_{t} f}$ is uniformly continuous. For this, let $x, y \in \mathbb{R}$. We have

$$
\begin{align*}
\frac{P_{t} f(x)}{1+x^{2}}-\frac{P_{t} f(y)}{1+y^{2}} & =\mathbb{E}\left[\frac{f\left(X_{t}^{0, x}\right)}{1+x^{2}}-\frac{f\left(X_{t}^{0, y}\right)}{1+y^{2}}\right]  \tag{A.4}\\
& =\mathbb{E}\left[I_{1}+I_{2}\right],
\end{align*}
$$

where

$$
\begin{aligned}
& I_{1}=\left(\tilde{f}\left(X_{t}^{0, x}\right)-\widetilde{f}\left(X_{t}^{0, y}\right)\right) \frac{1+\left(X_{t}^{0, x}\right)^{2}}{1+x^{2}} \\
& I_{2}=\widetilde{f}\left(X_{t}^{0, y}\right)\left(\frac{1+\left(X_{t}^{0, x}\right)^{2}}{1+x^{2}}-\frac{1+\left(X_{t}^{0, y}\right)^{2}}{1+y^{2}}\right) .
\end{aligned}
$$

Let $\epsilon>0$. By uniform continuity of $\tilde{f}$, there exists $\delta_{1}>0$ such that

$$
\begin{equation*}
\forall a, b \in \mathbb{R},|a-b| \leq \delta_{1} \Rightarrow|\widetilde{f}(a)-\widetilde{f}(b)|<\epsilon . \tag{A.5}
\end{equation*}
$$

Since $\lim _{M \rightarrow \infty} \mathbb{E}\left[\left|I_{1}\right| 1_{\left|\Gamma_{t}\right| \geq M}\right]=0$, there exists $M_{1}>0$ such that

$$
\begin{equation*}
\mathbb{E}\left[\left|I_{1}\right| 1_{\left|\Gamma_{t}\right| \geq M_{1}}\right]<\epsilon . \tag{A.6}
\end{equation*}
$$

We fix $0<\alpha<\frac{1}{2}$ and we choose $\delta_{2}=\left(\frac{\delta_{1}}{M_{1}}\right)^{1 / \alpha}$. Taking into account (A.1) and (A.5), for $|x-y|<\delta_{2}$ we have

$$
\begin{align*}
\mathbb{E}\left[\left|I_{1}\right| 1_{\left|\Gamma_{t}\right|<M_{1}}\right] & \leq \mathbb{E}\left[\frac{1+\left(X_{t}^{0, x}\right)^{2}}{1+x^{2}}\left(\widetilde{f}\left(X_{t}^{0, x}\right)-\widetilde{f}\left(X_{t}^{0, y}\right)\right) 1_{\left|X_{t}^{0, x}-X_{t}^{0, y}\right|<\delta_{1}}\right] \\
& <\sup _{\xi \in \mathbb{R}} \mathbb{E}\left[\frac{1+\left|X_{t}^{0, \xi}\right|^{2}}{1+\xi^{2}}\right] \epsilon . \tag{A.7}
\end{align*}
$$

The right-hand side is finite thanks to (A.2). Consequently, if $|x-y|<\delta_{2}$, then (A.6) implies that

$$
\begin{equation*}
\mathbb{E}\left[\left|I_{1}\right|\right]<A_{1} \epsilon, \tag{A.8}
\end{equation*}
$$

where $A_{1}=1+\sup _{\xi \in \mathbb{R}} \mathbb{E}\left[1+\frac{\left|X_{X}^{0, \xi}\right|^{2}}{1+\xi^{2}}\right]$.
Concerning $I_{2}$, we define

$$
\begin{equation*}
F(\omega, z)=\frac{1+\left|X_{t}^{0, z}(\omega)\right|^{2}}{1+z^{2}}, \omega \in \Omega, z \in \mathbb{R} \tag{A.9}
\end{equation*}
$$

Since $z \mapsto F(\cdot, z)$ is differentiable in $L^{2}(\Omega)$, by mean value theorem we get

$$
\begin{align*}
\mathbb{E}\left[\left|I_{2}\right|\right] & =|x-y| \mathbb{E}\left[\left|\widetilde{f}\left(X_{t}^{0, y}\right) \int_{0}^{1} \partial_{z} F(\cdot, a x+(1-a) y) d a\right|\right]  \tag{A.10}\\
& \leq|x-y|\|f\|_{E} \sup _{z} \mathbb{E}\left[\left|\partial_{z} F(\cdot, z)\right|\right] . \tag{A.11}
\end{align*}
$$

It remains to estimate the previous supremum. We have for $z \in \mathbb{R}$

$$
\partial_{z} F(\cdot, z)=2 \frac{X_{t}^{0, z} \partial_{z} X_{t}^{0, z}}{1+z^{2}}-2 z \frac{1+\left|X_{t}^{0, z}\right|^{2}}{\left(1+z^{2}\right)^{2}} .
$$

So by Cauchy-Schwarz we get

$$
\begin{align*}
\mathbb{E}\left[\left|\partial_{z} F(\cdot, z)\right|\right] & \leq 2\left(\frac{\mathbb{E}\left[\left|X_{t}^{0, z}\right|^{2}\right]}{1+z^{2}} \frac{\mathbb{E}\left[\left|\partial_{z} X_{t}^{0, z}\right|^{2}\right]}{1+z^{2}}\right)^{1 / 2}+2 \frac{|z|}{1+z^{2}} \frac{1+\mathbb{E}\left[\left|X_{t}^{0, z}\right|^{2}\right]}{1+z^{2}}  \tag{A.12}\\
& \leq A_{2} \tag{A.13}
\end{align*}
$$

where

$$
\begin{equation*}
A_{2}=2\left(\sup _{z} \frac{\mathbb{E}\left[\left|X_{t}^{0, z}\right|^{2}\right]}{1+z^{2}} \sup _{z} \mathbb{E}\left[\left|\partial_{z} X_{t}^{0, z}\right|^{2}\right]\right)^{1 / 2}+\left(1+\sup _{z} \frac{\mathbb{E}\left[\left|X_{t}^{0, z}\right|^{2}\right]}{1+z^{2}}\right) \tag{A.14}
\end{equation*}
$$

By (2.16) and (A.2) $A_{2}$ is finite and we get

$$
\begin{equation*}
\mathbb{E}\left[\left|I_{2}\right|\right] \leq A_{2}\|f\|_{E}|x-y| . \tag{A.15}
\end{equation*}
$$

Combining inequalities (A.8) and (A.15), (A.4) gives the existence of $\delta>0$ such that

$$
\begin{equation*}
|x-y|<\delta \Rightarrow\left|\frac{P_{t} f(x)}{1+x^{2}}-\frac{P_{t} f(y)}{1+y^{2}}\right|<\epsilon, \tag{A.16}
\end{equation*}
$$

so that the function $x \mapsto \frac{P_{t} f(x)}{1+x^{2}}$ is uniformly continuous.
In conclusion we have proved that $P_{t} f \in E . P_{t}$ is a bounded linear operator follows as a consequence of (A.3).

## Appendix B Proof of Theorem 2.18

We recall that the semigroup $P$ is here given by $P_{t} f(x)=\mathbb{E}\left[f\left(x+X_{t}\right)\right], x \in \mathbb{R}, t \geq 0$ and $X$ is a square integrable Lévy process vanishing at zero. The classical theory of semigroup for Lévy processes defines the semigroup $P$ on the set $C_{0}$ of continuous functions vanishing at infinity, equipped with the sup-norm $\|u\|_{\infty}=\sup _{x}|u(x)|$, c.f. for example [32, Theorem 31.5]. On $C_{0}$, the semigroup $P$ is strongly continuous, with norm $\|P\|=1$, and its generator $L_{0}$ is given by

$$
\begin{equation*}
L_{0} f(x)=\int\left(f(x+y)-f(x)-y f^{\prime}(x) 1_{|y|<1}\right) \nu(d y), f \in C_{0} . \tag{B.1}
\end{equation*}
$$

Moreover, [32, Theorem 31.5] shows that $C_{0}^{2} \subset D\left(L_{0}\right)$, where $C_{0}^{2}$ is the set of functions $f \in C^{2}$ such that $f, f^{\prime}$ and $f^{\prime \prime}$ vanish at infinity.

To prove Theorem 2.18 which concerns the infinitesimal generator of the semigroup $P$ defined on the set $E$ (c.f. (2.14)) related to a square integrable pure jump Lévy process, we adapt the classical theory. Since we consider a space $\left(E,\|\cdot\|_{E}\right)$, different from the classical one, i.e. $\left(C_{0},\|\cdot\|_{\infty}\right)$, we need to show that $\left(P_{t}\right)$ is still a strongly continuous semigroup.

Proposition B.1. Let $X$ be a square integrable Lévy process, then the semigroup $\left(P_{t}\right): E \rightarrow$ $E$ is strongly continuous.

Proof. The idea of the proof is an adaptation of the proof in [32, Theorem 31.5].
Let $f \in E$ and $\tilde{f}$ defined by $\tilde{f}(x)=\frac{f(x)}{1+x^{2}}, \forall x \in \mathbb{R}$. We evaluate, for $t>0, x \in \mathbb{R}$

$$
\frac{P_{t} f(x)-f(x)}{1+x^{2}}=\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right]+\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right) \frac{X_{t}^{2}+2 x X_{t}}{1+x^{2}}\right] .
$$

So

$$
\begin{equation*}
\left\|P_{t} f-f\right\|_{E} \leq \sup _{x \in \mathbb{R}}\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right]\right|+\sup _{x \in \mathbb{R}}\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right) \frac{X_{t}^{2}+2 x X_{t}}{1+x^{2}}\right]\right| . \tag{B.2}
\end{equation*}
$$

First, note that

$$
\begin{aligned}
\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right) \frac{X_{t}^{2}+2 x X_{t}}{1+x^{2}}\right]\right| & \leq\|f\|_{E} \mathbb{E}\left[\frac{X_{t}^{2}+2\left|x X_{t}\right|}{1+x^{2}}\right] \\
& \leq\|f\|_{E}\left(\mathbb{E}\left[X_{t}^{2}\right]+\mathbb{E}\left[\left|X_{t}\right|\right]\right)
\end{aligned}
$$

hence

$$
\sup _{x \in \mathbb{R}}\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right) \frac{X_{t}^{2}+2 x X_{t}}{1+x^{2}}\right]\right| \leq\|f\|_{E}\left(\mathbb{E}\left[X_{t}^{2}\right]+\mathbb{E}\left[\left|X_{t}\right|\right]\right) .
$$

Since $X$ is a square integrable Lévy process, $\mathbb{E}\left[X_{t}^{2}\right]=c_{2} t+c_{1}^{2} t^{2}$ where $c_{1}, c_{2}$ were defined in (2.25). Hence, the right-hand side of the inequality above goes to zero as $t$ goes to zero.

Now we prove that the first term $\sup _{x \in \mathbb{R}}\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right]\right|$ in the right-hand side of (B.2) goes to zero as well. Let $\epsilon>0$ be a fixed positive real. Since $\tilde{f}$ is uniformly continuous, then there is $\delta>0$ such that

$$
\forall x, y|x-y|<\delta \Rightarrow|\widetilde{f}(x)-\widetilde{f}(y)|<\frac{\epsilon}{2}
$$

Moreover, since $X$ is continuous in probability

$$
\exists t_{0}>0, \text { such that } \forall t<t_{0}, \mathbb{P}\left(\left|X_{t}\right|>\delta\right)<\frac{\epsilon}{4\|f\|_{E}}
$$

For all $x \in \mathbb{R}, t<t_{0}$ we have

$$
\begin{aligned}
\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right]\right| & \leq \mathbb{E}\left[\left|\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right| 1_{\left\{\left|X_{t}\right| \leq \delta\right\}}\right]+\mathbb{E}\left[\left|\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right| 1_{\left\{\left|X_{t}\right|>\delta\right\}}\right] \\
& \leq \frac{\epsilon}{2}+2\|f\|_{E} \mathbb{P}\left(\left|X_{t}\right|>\delta\right) \\
& \leq \epsilon
\end{aligned}
$$

Since the inequality above is valid for every $x \in \mathbb{R}$, then

$$
\sup _{x \in \mathbb{R}}\left|\mathbb{E}\left[\widetilde{f}\left(x+X_{t}\right)-\widetilde{f}(x)\right]\right| \xrightarrow{t \rightarrow 0} 0
$$

This concludes the proof that $P$ is a strongly continuous semigroup.
Remark B.2. Note that the semigroup $\left(P_{t}\right)$ is not a contraction. In fact, if $f \in E, t>0$, then

$$
\begin{equation*}
\left\|P_{t} f\right\|_{E}=\sup _{x \in \mathbb{R}}\left|\mathbb{E}\left[\frac{f\left(x+X_{t}\right)}{1+x^{2}}\right]\right| \tag{B.3}
\end{equation*}
$$

Let $f_{0}(x)=1+x^{2}$ and denote again $c_{1}=\mathbb{E}\left[X_{1}\right]$ and $c_{2}=\operatorname{Var}\left(X_{1}\right)$. Obviously $f_{0} \in E$, $\left\|f_{0}\right\|_{E}=1$ and

$$
\begin{align*}
\left\|P_{t} f_{0}\right\|_{E} & =\sup _{x \in \mathbb{R}} \mathbb{E}\left[\frac{1+\left(x+X_{t}\right)^{2}}{1+x^{2}}\right] \\
& =1+\sup _{x \geq 0} \frac{2 x\left|c_{1}\right| t+c_{2} t+c_{1}^{2} t^{2}}{1+x^{2}}  \tag{B.4}\\
& =1+\left|c_{1}\right| t+c_{1}^{2} t^{2}
\end{align*}
$$

Hence $\left(P_{t}\right)$ cannot be not a contraction since

$$
\left\|P_{t}\right\| \geq\left\|P_{t} f_{0}\right\|_{E}>1
$$

On the other hand, for $f \in E$, (B.3) gives

$$
\left\|P_{t} f\right\|_{E} \leq\|f\|_{E}\left\|P_{t} f_{0}\right\|_{E}
$$

By (B.4) this implies that

$$
\left\|P_{t}\right\| \leq 1+\left(\left|c_{1}\right|+c_{2}\right) t+c_{1}^{2} t^{2}
$$

So, there exists a positive real $\omega>0$ such that

$$
\left\|P_{t}\right\| \leq e^{\omega t}
$$

Semigroups verifying the latter inequality are called quasi-contractions, see [29]. For instance, [29, Corollary 3.8] implies that

$$
\begin{equation*}
\forall \lambda>\omega, \lambda I-L \text { is invertible. } \tag{B.5}
\end{equation*}
$$

At this point we show that the space $E_{0}^{2}$, defined in (2.23), is a subset of $D(L)$ and that formula (B.1) remains valid in $E_{0}^{2}$. This will be done adapting a technique described in [32, Theorem 31.5], where it is stated that $C_{0}^{2}$ is included in $D\left(L_{0}\right)$. The main tool used for the proof of [32, Theorem 31.5] is the small time asymptotics

$$
\begin{equation*}
\lim _{t \rightarrow 0} \frac{1}{t} \mathbb{E}\left[g\left(X_{t}\right)\right]=\int g(x) \nu(d x) \tag{B.6}
\end{equation*}
$$

which holds for bounded continuous function $g$ vanishing on a neighborhood of the origin, see [32, Corollary 8.9]. This result has been extended to a class of unbounded functions by [12, Theorem 1.1]. (B.6) is used in [12, Proposition 2.3] to prove that the quantity $\lim _{t \rightarrow 0} \frac{P_{t} g-g}{t}(x)$ converges point-wise, under some suitable conditions on the function $g$.

We state a similar lemma below.
Lemma B.3. Let $f \in E_{0}^{2}$. For all $x \in \mathbb{R}$, the quantity

$$
\begin{equation*}
\lim _{t \rightarrow 0} \frac{P_{t} f-f}{t}(x) \tag{B.7}
\end{equation*}
$$

exists and equals the right-hand side of (B.1).

## Remark B.4.

1) To be self-contained, we give below a simple proof of Lemma B.3, in the case when $X$ is a square integrable pure jump process.
2) Later we will need to show that the point-wise convergence (B.7) holds according to the norm of $E$.

Proof. Let $f \in E_{0}^{2}$. First, we verify that the integral

$$
\begin{equation*}
\int\left(f(x+y)-f(x)-y f^{\prime}(x) 1_{|y|<1}\right) \nu(d y) \tag{B.8}
\end{equation*}
$$

is well-defined for all $x \in \mathbb{R}$, taking into account $\int y^{2} \nu(d y)<\infty$ by (2.24).
In fact, by Taylor expansion and since $f \in E_{0}^{2}$, then for every $x \in \mathbb{R}$ there exist $a, b \geq 0$ such that, for all $y \in \mathbb{R}$

$$
\begin{aligned}
|f(x+y)-f(x)| 1_{|y| \geq 1} & \leq a\left(y^{2}+1\right) 1_{|y| \geq 1} \\
\left|f(x+y)-f(x)-f^{\prime}(x) y\right| 1_{|y|<1} & \leq b y^{2} 1_{|y|<1}
\end{aligned}
$$

Let $t>0, x \in \mathbb{R}$. By Taylor expansion and Fubini theorem, recalling that $P_{t} f(x)=$ $\mathbb{E}\left[f\left(x+X_{t}\right)\right]$ we have

$$
\frac{P_{t} f-f}{t}(x)=c_{1} f^{\prime}(x)+\int_{0}^{1}(1-a) \frac{1}{t} \mathbb{E}\left[f^{\prime \prime}\left(a X_{t}+x\right) X_{t}^{2}\right] d a
$$

By abuse of notation, we denote by $L_{0} f(x)$ the integral (B.8). Taking into account (2.25) we have

$$
\begin{align*}
L_{0} f(x) & =c_{1} f^{\prime}(x)+\int\left(f(x+y)-f(x)-y f^{\prime}(x)\right) \nu(d y) \\
& =c_{1} f^{\prime}(x)+\int_{0}^{1}(1-a) \int_{\mathbb{R}} y^{2} f^{\prime \prime}(a y+x) \nu(d y) d a \tag{B.9}
\end{align*}
$$

Hence, it remains to show that ( $x$ being fixed)

$$
\begin{equation*}
\frac{P_{t} f-f}{t}(x)-L_{0} f(x)=\int_{0}^{1}(1-a)\left(\frac{1}{t} \mathbb{E}\left[X_{t}^{2} f^{\prime \prime}\left(a X_{t}+x\right)\right]-\int_{\mathbb{R}} y^{2} f^{\prime \prime}(a y+x) \nu(d y)\right) d a \tag{B.10}
\end{equation*}
$$

$$
\xrightarrow[t \rightarrow 0]{ } 0
$$

For $a \in[0,1]$, we denote $g(y)=y^{2} f^{\prime \prime}(a y+x)$. We have $g(y) \underset{y \rightarrow 0}{\sim} y^{2} f^{\prime \prime}(x)$. If $f^{\prime \prime}(x) \neq 0$, then [12, Theorem 1.1] (ii) implies that

$$
\begin{equation*}
\lim _{t \rightarrow 0} \frac{1}{t} \mathbb{E}\left[g\left(X_{t}\right)\right]=\int_{\mathbb{R}} g(y) \nu(d y) \tag{B.11}
\end{equation*}
$$

If $f^{\prime \prime}(x)=0$, then $g(y)=o\left(y^{2}\right)$ and (B.11) is still valid by [12, Theorem 1.1] (i). We conclude to the validity of (B.10) by Lebesgue dominated convergence theorem taking into account that $f^{\prime \prime}$ is bounded.

As observed in a similar case in [12, Remark 2.4], we will prove that the point-wise convergence proved in Lemma B. 7 holds in the strong sense.

For this purpose, we introduce the linear subspace

$$
\begin{equation*}
\widetilde{E}=\left\{f \in C \text { such that } \widetilde{f}:=x \mapsto \frac{f(x)}{1+x^{2}} \text { is vanishing at infinity }\right\} \tag{B.12}
\end{equation*}
$$

of $E$. It it is easy to show that $\widetilde{E}$ is closed in $E$ so that it is a Banach subspace of $E$.
Lemma B.5. Let $f, g \in \widetilde{E}$, such that

$$
\begin{equation*}
\lim _{t \rightarrow 0} \frac{P_{t} f-f}{t}(x)=g(x), \forall x \in \mathbb{R} \tag{B.13}
\end{equation*}
$$

Then $f \in D(L)$ and $L f=g$.
Proof. We first introduce a restriction $\widetilde{P}$ of the semigroup $P$ to the linear subspace $\widetilde{E}$. By Lebesgue dominated convergence theorem and the fact that $\frac{1+\left(X_{t}+x\right)^{2}}{1+x^{2}} \leq 2\left(\left|X_{t}\right|^{2}+1\right)$, one can show that $P_{t} f \in \widetilde{E}$ for any $f \in \widetilde{E}, t \geq 0$. Hence $\left(\widetilde{P}_{t}\right)$ is a semigroup on $\widetilde{E}$; we denote by $\widetilde{L}$ its infinitesimal generator.

As in [32, Lemma 31.7], we denote by $L^{\#} f=g$, the operator defined by the equation (B.13) for $f, g \in \widetilde{E}$ and by $D\left(L^{\#}\right)$ its domain, i.e. the set of functions $f$ for which (B.13) exists. Then $L^{\#}$ is an extension of $\widetilde{L}$.

Fix $q>\left|c_{1}\right|+c_{2}$. We prove first that

$$
\begin{equation*}
\forall f \in D\left(L^{\#}\right) \quad\left(q I-L^{\#}\right) f=0 \Rightarrow f=0 \tag{B.14}
\end{equation*}
$$

Let $f \in D\left(L^{\#}\right)$ such that $\left(q I-L^{\#}\right) f=0$. We denote $f^{-}=-(f \wedge 0)$ and $f^{+}=f \vee 0$. Suppose that $f^{+} \neq 0$. Since $\widetilde{f^{+}}$is continuous and vanishing at infinity, there exists $x_{1}$ such that $\frac{f^{+}\left(x_{1}\right)}{1+x_{1}^{2}}=\max _{x} \frac{f^{+}(x)}{1+x^{2}}>0$. Moreover $f\left(x_{1}\right)=f^{+}\left(x_{1}\right)$. Then

$$
\frac{\mathbb{E}\left[f\left(x_{1}+X_{t}\right)-f\left(x_{1}\right)\right]}{t} \leq \frac{1}{t}\left(f\left(x_{1}\right) \frac{\mathbb{E}\left[1+\left(x_{1}+X_{t}\right)^{2}\right]}{1+x_{1}^{2}}-f\left(x_{1}\right)\right) .
$$

Passing to the limit when $t \rightarrow 0$ it follows

$$
L^{\#} f\left(x_{1}\right) \leq f\left(x_{1}\right)\left(\left|c_{1}\right|+c_{2}\right) .
$$

Then

$$
\left(q-\left|c_{1}\right|-c_{2}\right) f\left(x_{1}\right) \leq 0,
$$

which contradicts the fact that $f\left(x_{1}\right)>0$. Hence, $f^{+}=0$. With similar arguments, we can show that $f^{-}=0$ and so $f=0$, which proves (B.14).

By restriction, $\left(\widetilde{P}_{t}\right)$ fulfills $\left\|\widetilde{P}_{t}\right\| \leq e^{\omega t}$, in particular it is a quasi-contraction semigroup, so by (B.5), we can certainly choose $q>\max \left(\left|c_{1}\right|+c_{2}, \omega\right)$, so that $q I-\widetilde{L}$ is invertible and $R(q I-\widetilde{L})=\widetilde{E}$.
We observe that $D(\widetilde{L}) \subset D\left(L^{\#}\right)$. Let now $f \in D\left(L^{\#}\right)$; then $\left(q I-L^{\#}\right) f \in \widetilde{E}=R(q I-\widetilde{L})$. Consequently, there is $v \in D(\widetilde{L})$ such that $\left(q I-L^{\#}\right) f=(q I-\widetilde{L}) v$. So, $\left(q I-L^{\#}\right)(f-v)=0$.

By (B.14), $\left(q I-L^{\#}\right)$ is injective, so $f=v$ and $f \in D(\widetilde{L})$. Consequently $\widetilde{L} f$ is given by $g$ defined in (B.13). Finally, the fact that $D(\widetilde{L}) \subset D(L)$ and $\widetilde{L}$ is a restriction of $L$ allow to conclude the proof of Lemma B.5.

We continue the proof of Theorem 2.18 making use of Lemmas B. 3 and Lemma B.5. First, let us prove that $E_{0}^{2} \subset \widetilde{E}$. Indeed by Taylor expansion, we have, for $f \in E_{0}^{2}$

$$
\frac{f(x)}{1+x^{2}}=\frac{f(0)}{1+x^{2}}+\frac{x}{1+x^{2}} f^{\prime}(0)+\frac{x^{2}}{1+x^{2}} \int_{0}^{1}(1-\alpha) f^{\prime \prime}(x \alpha) d \alpha
$$

Since $\lim _{x \rightarrow \infty} f^{\prime \prime}(x \alpha)=0$ for all $\left.\alpha \in\right] 0,1[$, then by Lebesgue theorem, we have that $\lim _{x \rightarrow \infty} \frac{f(x)}{1+x^{2}}=0$, so $f \in \widetilde{E}$.

By Lemma B.3, it follows

$$
\lim _{t \rightarrow 0} \frac{P_{t} f-f}{t}(x)=L_{0} f(x), \forall x \in \mathbb{R}
$$

where $L_{0}$ is given in (B.8). In order to apply Lemma B.5, it remains to show that $L_{0} f \in \widetilde{E}$. Using relation (B.9), for $x \in \mathbb{R}$ we get

$$
\frac{L_{0} f(x)}{1+x^{2}}=c_{1} \frac{f^{\prime}(x)}{1+x^{2}}+\int_{0}^{1}(1-a) \int_{\mathbb{R}} y^{2} \frac{f^{\prime \prime}(a y+x)}{1+x^{2}} \nu(d y) d a .
$$

Since $f \in E_{0}^{2}$, then $f^{\prime \prime}$ is bounded and $f^{\prime}$ has linear growth. So, the fact that $\int_{\mathbb{R}} y^{2} \nu(d y)<\infty$ implies indeed $\lim _{x \rightarrow \infty} \frac{L_{0} f(x)}{1+x^{2}}=0$ and $L_{0} f \in \widetilde{E}$.

Finally, Lemma B. 5 implies that $E_{0}^{2} \subset D(L)$ and for $f \in E_{0}^{2}, L f$ is given by (2.27).

ACKNOWLEDGEMENTS: Financial support was provided by the ANR Project MASTERIE 2010 BLAN-0121-01. The authors are also grateful to Prof. Robert Dalang for the support.

## References

[1] S. Ankirchner and G. Heyne. Cross hedging with stochastic correlation. Finance Stoch., 16(1):17-43, 2012.
[2] M. Barlow and M. Yor. Semi-martingale inequalities via the Garsia-Rodemich-Rumsey lemma, and applications to local times. Journal of Functional Analysis, 49(2):198 - 229, 1982.
[3] J.-M. Bismut. Conjugate convex functions in optimal stochastic control. J. Math. Anal. Appl., 44:384-404, 1973.
[4] P. Briand, B. Delyon, and J. Mémin. On the robustness of backward stochastic differential equations. Stochastic Process. Appl., 97(2):229-253, 2002.
[5] R. Buckdahn. Backward stochastic differential equations driven by a martingale. Unpublished, 1993.
[6] R. Carbone, B. Ferrario, and M. Santacroce. Backward stochastic differential equations driven by càdlàg martingales. Teor. Veroyatn. Primen., 52(2):375-385, 2007.
[7] A. Černý and J. Kallsen. On the structure of general mean-variance hedging strategies. Ann. Probab., 35(4):1479-1531, 2007.
[8] R. Cont and P. Tankov. Financial modelling with jump processes. Chapman \& Hall/CRC Financial Mathematics Series. Chapman \& Hall/CRC, Boca Raton, FL, 2004.
[9] M. H. A. Davis. Option hedging with basis risk. From Stochastic Calculus to Mathematical Finance, Springer-Verlag, 2006.
[10] N. El Karoui, S. Hamadène, and A. Matoussi. Backward stochastic differential equations and applications. In R. Carmona, editor, Indifference pricing: theory and applications, chapter 8, pages 267-320. Springer-Verlag, 2008.
[11] N. El Karoui and S.-J. Huang. A general result of existence and uniqueness of backward stochastic differential equations. In Backward stochastic differential equations (Paris,

1995-1996), volume 364 of Pitman Res. Notes Math. Ser., pages 27-36. Longman, Harlow, 1997.
[12] J. E. Figueroa-López. Small-time moment asymptotics for Lévy processes. Statist. Probab. Lett., 78(18):3355-3365, 2008.
[13] H. Föllmer and M. Schweizer. Hedging of contingent claims under incomplete information. In Applied stochastic analysis (London, 1989), volume 5 of Stochastics Monogr., pages 389-414. Gordon and Breach, New York, 1991.
[14] A. Friedman. Partial differential equations of parabolic type. Prentice-Hall, Inc., Englewood Cliffs, N.J., 1964.
[15] A. Friedman. Stochastic differential equations and applications. Vol. 1. Academic Press [Harcourt Brace Jovanovich, Publishers], New York-London, 1975. Probability and Mathematical Statistics, Vol. 28.
[16] S. Goutte, N. Oudjane, and F. Russo. On some expectation and derivative operators related to integral representations of random variables with respect to a PII process. Stoch. Anal. Appl., 31(1):108-141, 2013.
[17] S. Goutte, N. Oudjane, and F. Russo. Variance optimal hedging for continuous time additive processes and applications. Stochastics, 86(1):147-185, 2014.
[18] V. Henderson and D. Hobson. Substitute hedging. Risk, 15:171-75, 2002.
[19] F. Hubalek, J. Kallsen, and L. Krawczyk. Variance-optimal hedging for processes with stationary independent increments. Ann. Appl. Probab., 16(2):853-885, 2006.
[20] H. Hulley and T. A. McWalter. Quadratic hedging of basis risk. Research Paper Series 225, Quantitative Finance Research Centre, University of Technology, Sydney, 2008.
[21] N. Jacob. Pseudo differential operators and Markov processes. Vol.I. Imperial College Press, London, 2001. Fourier analysis and semigroups.
[22] J. Jacod and A. N. Shiryaev. Limit theorems for stochastic processes, volume 288 of Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. Springer-Verlag, Berlin, second edition, 2003.
[23] I. Karatzas and S. E. Shreve. Brownian motion and stochastic calculus, volume 113 of Graduate Texts in Mathematics. Springer-Verlag, New York, second edition, 1991.
[24] P. Monat and C. Stricker. Föllmer-Schweizer decomposition and mean-variance hedging for general claims. Ann. Probab., 23(2):605-628, 1995.
[25] M. Monoyios. Performance of utility-based strategies for hedging basis risk. Quant. Finance, 4(3):245-255, 2004.
[26] M. Monoyios. Optimal hedging and parameter uncertainty. IMA J. Manag. Math., 18(4):331-351, 2007.
[27] É. Pardoux and S. G. Peng. Adapted solution of a backward stochastic differential equation. Systems Control Lett., 14(1):55-61, 1990.
[28] É. Pardoux and A. Rascanu. Stochastic differential equations, Backward SDEs, Partial differentialequations. Springer, 2014.
[29] A. Pazy. Semigroups of linear operators and applications to partial differential equations, volume 44 of Applied Mathematical Sciences. Springer-Verlag, New York, 1983.
[30] P. E. Protter. Stochastic integration and differential equations, volume 21 of Stochastic Modelling and Applied Probability. Springer-Verlag, Berlin, 2005. Second edition. Version 2.1, Corrected third printing.
[31] F. Russo and G. Trutnau. Some parabolic PDEs whose drift is an irregular random noise in space. Ann. Probab., 35(6):2213-2262, 2007.
[32] K.-i. Sato. Lévy processes and infinitely divisible distributions, volume 68 of Cambridge Studies in Advanced Mathematics. Cambridge University Press, Cambridge, 2013. Translated from the 1990 Japanese original, Revised edition of the 1999 English translation.
[33] M. Schweizer. Approximating random variables by stochastic integrals. Ann. Probab., 22(3):1536-1575, 1994.
[34] M. Schweizer. A guided tour through quadratic hedging approaches. In Option pricing, interest rates and risk management, Handb. Math. Finance, pages 538-574. Cambridge Univ. Press, Cambridge, 2001.
[35] D. W. Stroock and S. Varadhan. Multidimensional diffusion processes. Classics in Mathematics. Springer-Verlag, Berlin, 2006. Reprint of the 1997 edition.

