# When Capital Is a Funding Source: The Anticipated Backward Stochastic Differential Equations of X-Value Adjustments 

Stéphane Crépey, Wissal Sabbagh, Shiqi Song

## - To cite this version:

Stéphane Crépey, Wissal Sabbagh, Shiqi Song. When Capital Is a Funding Source: The Anticipated Backward Stochastic Differential Equations of X-Value Adjustments. SIAM Journal on Financial Mathematics, 2020, 11, pp.99-130. 10.1137/19m1242781 . hal-03910119

HAL Id: hal-03910119

## https://hal.science/hal-03910119

Submitted on 21 Dec 2022

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.

# When Capital Is a Funding Source: The Anticipated Backward Stochastic Differential Equations of X-Value Adjustments* 

Stéphane Crépey ${ }^{\dagger}$, Wissal Sabbagh ${ }^{\dagger}$, and Shiqi Song ${ }^{\dagger}$


#### Abstract

X-value adjustments (XVAs) refer to various financial derivative pricing adjustments accounting for counterparty risk and its funding (FVA) and capital (KVA) implications for a bank. In this paper we show that the XVA equations are well-posed, including in the realistic case where capital is deemed fungible as a source of funding for variation margin. This intertwining of capital at risk and the FVA, added to the fact that the KVA is part of capital at risk, leads to a system of backward SDEs (BSDEs) of the McKean type (anticipated BSDEs) for the FVA and the KVA, with coefficients entailing a conditional risk measure of the one-year-ahead increment of the martingale part of the FVA. This is first considered in the case of a hypothetical bank without debt. In the practical case of a defaultable bank, the resulting anticipated BSDEs, which are stopped before the default of the bank, are solved likewise after reduction to a reference market filtration.


Key words. credit valuation adjustment (CVA), funding valuation adjustment (FVA), capital valuation adjustment (KVA), anticipated (or McKean) BSDE, progressive enlargement of filtration, invariance time

AMS subject classifications. 91B30, 91G20, 91G40, 60G07, 60G44
DOI. 10.1137/19M1242781

1. Introduction. XVAs, where VA stands for valuation adjustment and X is a catch-all letter to be replaced by C for credit, D for debt, F for funding, M for margin, or K for capital, denote various pricing adjustments applied to financial derivatives since the 2008 crisis, in order to account for counterparty risk and its capital and funding implications for a bank.

The valuation of securities in a defaultable environment has started with the pricing of risky bonds or credit default swaps (CDS). These only involve a single default time, the one of the issuer of the bond or of the reference firm in a CDS; see Pykhtin (2005) for a collection of early CVA papers. But the 2008 crisis clearly showed that banks are themselves risky. As a consequence, new regulations have been established on the functioning of financial institutions (see Basel Committee on Banking Supervision (2011)), with a great impact on capital, margin, and clearing requirements for investment banks, making it vital for them to understand how they should operate in this new environment.

[^0]The question has first been studied in a setup where the default risk of the bank and a counterparty are treated symmetrically, in a CVA/DVA mindset (see, e.g., Brigo and Capponi (2010) and cf. previously Duffie and Huang (1996) or Bielecki and Rutkowski (2002, (14.25) p. 448)). However, the view of the industry evolved quite dramatically, banks reacting to the Basel III regulatory changes by pricing further the FVA, the MVA, and the KVA, which are the respective costs of funding the so called variation margin (VM), initial margin (IM), and of capital; see Brigo and Pallavicini (2014), Bichuch, Capponi, and Sturm (2018), and Burgard and Kjaer (2011, 2013, 2017) (without KVA) or, with a KVA meant as a liability like the CVA and the FVA (as opposed to a risk premium in our case), Green, Kenyon, and Dennis (2014) and Elouerkhaoui (2016). See also Andersen, Duffie, and Song (2019) and Bielecki and Rutkowski (2015) for different focuses on the funding side of the problem, with respective emphases on the related wealth transfer and nonlinear arbitrage issues.

The main dividing line in this literature is between an XVA replication approach and a cost-of-capital, incomplete market approach. Following up on the Hull and White (2012a), Hull and White (2012b) prompted FVA debate, Albanese and Andersen (2015), Albanese, Caenazzo, and Crépey (2016), and Albanese and Crépey (2019) have delineated the specific implications of the default risk of the bank itself, providing a better insight on counterparty risk considered not only in the strict sense, but also through its consequences in terms of capital and funding. A key point in Albanese and Crépey (2019) is that, in order to account for the defaultability of the bank itself, all the cumulative cash flow and value processes must be stopped before the bank default time $\tau$ in the XVA equations. Indeed, given the impossibility for bank shareholders to hedge these cash flows and monetize them before $\tau$, later cash flows only benefit bank bondholders.

In Albanese and Crépey (2019, section 5, Remark 4.3), the XVA equations are studied in the base case where capital at risk is not used by the bank for its funding purposes. Then, in the context of partially or uncollateralized transactions, the FVA can seem very large. For instance, in January 2014, JP Morgan has recorded a $\$ 1.5$ billion FVA loss. ${ }^{1}$ However, in practice, capital at risk can be used for funding the so called VM (cf. the parameter $\phi$ representing "the fraction of capital used for funding" in Green, Kenyon, and Dennis (2014)). This may cause a material FVA reduction, as high as one-half or more on a real banking portfolio, as demonstrated numerically in Albanese, Caenazzo, and Crépey (2017, section 5.2) (see also Albanese et al. (2019, section 5.2)).

In this paper, we provide a mathematical analysis of the FVA and its KVA implications when capital at risk is a possible funding source. Our results thus complement the XVA analysis of Albanese and Crépey (2019, section 5), which had been made under the assumption that capital at risk is not used for funding purposes, as well as the numerical studies of Albanese, Caenazzo, and Crépey (2017, section 5.2) and Albanese et al. (2019, section 5) which are based on the equations of the present paper.

The fungibility of capital at risk as a source of funding for VM leads to a system of anticipated backward stochastic differential equations (ABSDEs, or BSDEs "of the McKean type"). Peng and Yang (2009) have introduced ABSDEs in relation to a problem of SDEs with delay. They established the well-posedness of a multivariate ABSDE on a fixed time horizon with Lipschitz coefficients in a Brownian setup. As usual, in the univariate case, the Lipschitz

[^1]condition can be relaxed into continuous coefficients with linear growth, which was done in Elliot and Yang (2013). In parallel to our work, an ABSDE involving a conditional expected shortfall as anticipated term (by contrast with a conditional expectation in the previous ABSDE literature) has been considered in Agarwal et al. (2019) to study the cost of IM, in a univariate, Brownian, and Lipschitz setup, over a fixed time horizon. In this paper, we solve a system of ABSDEs with jumps (because of the counterparty defaults), monotone coefficients (this is in fact for the sake of generality; our XVA coefficients are even Lipschitz), a more general anticipated dependence of the coefficient (on the integrand components of the solution), and stopped before a random time (the default time of the bank itself, in the XVA context).

The outline of the paper is as follows. Section 2 revisits the XVA equations from Albanese and Crépey (2019) when capital is deemed fungible as a source of funding for VM, first assuming the bank default free. Section 3 establishes a general ABSDE well-posedness result, which is applied in section 4 to show the well-posedness of the XVA ABSDEs in the theoretical case of a default-free bank. Section 5 , which is of independent interest, extends beyond the basic immersion setup the classical credit risk intensity pricing formulas (see, e.g., Bielecki, Jeanblanc, and Rutkowski (2009, Chapter 3)). This framework is used in section 6 for stating and solving the XVA ABSDEs in the realistic case of a defaultable bank. Section 7 wraps up the paper.

The main contributions of the paper are Theorems 3.1 (an extension with jumps, monotone coefficient, and a more general anticipated dependence, of the ABSDE result of Peng and Yang (2009)), 4.1 (well-posedness of the XVA ABSDEs in the theoretical case of a default-free bank), 5.1 (extension to the invariance time setup of the classical credit risk intensity pricing formulas), and 6.1 (well-posedness of the XVA ABSDEs in the realistic case of a defaultable bank).

### 1.1. Mathematical setup. We denote by

- . ${ }^{\top}$, vector or matrix transpose;
- $|\cdot|$ and $\langle\cdot, \cdot\rangle$, Euclidean norms and scalar products in the dimension of their arguments (vectors or vectorized matrices);
- $(\Omega, \mathfrak{A}, \mathbb{F}, \mathbb{P})$, a filtered probability space, for a complete and right continuous filtration $\mathbb{F}=\left(\mathfrak{F}_{t}\right)_{t \geq 0}$ of sub- $\sigma$ fields of a reference $\sigma$ field $\mathfrak{A}$ and for a probability measure $\mathbb{P}$;
- $\mathbb{E}$ and $\mathbb{E}_{t}$, the $\mathbb{P}$ expectation and the $\left(\mathfrak{F}_{t}, \mathbb{P}\right)$ conditional expectation;
- $\mathfrak{P}$, the $\mathbb{F}$ predictable sigma-field;
- $\mathfrak{B}(E)$, the Borel $\sigma$ field on any metrizable space $E$;
- $m(S)$, the $(\mathbb{F}, \mathbb{P})$ canonical Doob-Meyer local martingale component of an $(\mathbb{F}, \mathbb{P})$ special semimartingale $S$ (with $\left.m\left(S_{0}\right)=S_{0}\right)$;
- $C$, a positive constant, the value of which may change from line to line.

Stochastic integrals of $\mathfrak{P}$ measurable ( $\mathbb{F}$ predictable) processes against semimartingales and stochastic integrals of $\widehat{\mathfrak{P}}$ measurable random functions with respect to jump measures or their compensations are defined as in Jacod (1979). Stochastic integrals are sometimes written in - notation, using the precedence convention $K L \cdot X=(K L)$. $X$.

As can be classically established by section theorem, for any progressive process (Lebesgue integrand) $X$ such that the predictable projection ${ }^{p} X$ exists, ${ }^{2}$ the indistinguishable equality $\int_{0}^{\cdot} p X_{s} d s=\int_{0}^{*} X_{s} d s$ holds. As a consequence, we only consider predictable Lebesgue integrands (even if this means replacing $X$ by ${ }^{p} X$ ).

[^2]We denote by $B$ an $(\mathbb{F}, \mathbb{P})$ standard $d$ variate Brownian motion for some nonnegative integer $d$. We denote by $E$ an Euclidean space, by $\pi$ a $\sigma$ finite measure on $(E, \mathfrak{B}(E))$ such that $\int_{E}\left(1 \wedge|e|^{2}\right) \pi(d e)<\infty$, and we write $\widehat{\mathfrak{P}}=\mathfrak{P} \otimes \mathfrak{B}(E)$. We consider an $\mathbb{F}$ optional integer valued random measure $j(d t, d e)$ on $\mathbb{R}_{+} \times E$, with $\mathbb{P}$ compensator $\eta(t, e) \pi(d e) d t$ and compensated martingale measure $M(d t, d e)$, for some nonnegative and bounded $\mathfrak{P}$ measurable random function $\eta$.

Given a positive integer $l$, we introduce

- $\mathcal{S}_{2}^{l}$, the space of $\mathbb{R}^{l}$ valued $\mathbb{F}$ adapted càdlàg processes $Y$ such that

$$
\begin{equation*}
\|Y\|_{\mathcal{S}_{2}^{l}}^{2}=\mathbb{E}\left[\sup _{0 \leq t \leq T}\left|Y_{t}\right|^{2}\right]<+\infty \tag{1}
\end{equation*}
$$

and $\mathcal{M}_{2}^{l}$, the space of martingales (componentwise) in $\mathcal{S}_{2}^{l}$;

- $\mathcal{H}_{2}^{l}$, the space of $\mathbb{R}^{l \otimes d}$ valued $\mathbb{F}$ predictable processes $Z$ such that

$$
\|Z\|_{\mathcal{H}_{2}^{l}}^{2}=\mathbb{E}\left[\int_{0}^{T}\left|Z_{t}\right|^{2} d t\right]<+\infty
$$

- $\mathcal{L}_{0}^{l}$, the space of $l$-variate $\mathfrak{B}(E)$ measurable functions endowed with the topology of convergence in measure induced by $\pi$ and we write, for any time $t$ and $u \in \mathcal{L}_{0}^{l}$,

$$
|u|_{t}=\left(\int_{E}|u(e)|^{2} \eta(t, e) \pi(d e)\right)^{1 / 2}
$$

- $\widehat{\mathcal{H}}_{2}^{l}$, the space of $l$-variate $\widehat{\mathfrak{P}}$ measurable random functions $U$ such that

$$
\|U\|_{\hat{\mathcal{H}}_{2}^{l}}^{2}=\mathbb{E}\left[\int_{0}^{T} \int_{E}\left|U_{t}(e)\right|^{2} \eta(t, e) \pi(d e) d t\right]=\mathbb{E}\left[\int_{0}^{T}\left|U_{t}\right|_{t}^{2} d t\right]<\infty .
$$

In the case where $l=1$ we drop the index $l$, e.g., we write $\mathcal{S}_{2}$ instead of $\mathcal{S}_{2}^{1}$. We introduce likewise the space $\mathcal{H}_{1}$ of real valued $\mathbb{F}$ predictable processes $X$ such that

$$
\begin{equation*}
\|X\|_{\mathcal{H}_{1}}=\mathbb{E}\left[\int_{0}^{T}\left|X_{t}\right| d t\right]<+\infty \tag{2}
\end{equation*}
$$

We assume that every $(\mathbb{F}, \mathbb{P})$ square integrable martingale null at time 0 has a representation of the form

$$
\begin{equation*}
\int_{0}^{t} Z_{s} d B_{s}+\int_{0}^{t} \int_{E} U_{s}(e) M(d s, d e), 0 \leq t \leq T \tag{3}
\end{equation*}
$$

for suitable integrands $Z \in \mathcal{H}_{2}$ and $U \in \widehat{\mathcal{H}}_{2}$, uniquely defined modulo $d[B, B]$ and $\eta(t, e) \pi(d e) d t$ negligible sets, respectively. The left and right terms in (3) represent the corresponding vector and random measure stochastic integrals.

Given this square integrable martingale representation assumption, one can readily check that all the results in Kruse and Popier $(2016 \text {, section } 4)^{3}$ derived there in the case of a Poisson random measure, are valid in our setup.

[^3]2. The XVA equations in the case of a default-free bank. Let there be given

- $r$, an $\mathbb{F}$ progressive risk-free interest rate process, with related risk-free discount factor $\beta=e^{-\int_{0}^{j} r_{t} d t}$;
- $\lambda$, an $\mathbb{F}$ progressive bank funding (borrowing) spread process,
both assumed bounded (possibly after truncation by a large positive constant of standard interest rate or credit spread models).

We consider a derivative portfolio with final maturity $T$ between a bank and a client. In order to focus on counterparty risk and XVA analysis, we assume that the bank has set up a perfectly collateralized market hedge of its client portfolio (see section 3.2 in Albanese and Crépey (2019) for a more detailed description). Hence only the counterparty risk related cash flows remain.

In this section, we assume that only the client is default prone (the defaultability of the bank itself will be added in section 6). The probability measure $\mathbb{P}$ is then interpreted as a riskneutral pricing measure, calibrated to prices of liquid derivatives as standard. In the context of a cost-of-capital XVA approach, the historical probability measure required for capital at risk computations is then taken equal to $\mathbb{P}$ (see Remark 3.5 in Albanese and Crépey (2019)), leaving the discrepancy between the two to model risk.

Remark 2.1. In the hypothetical case of a default-free bank, the bank borrowing spread $\lambda$ is interpreted as a liquidity spread. However, in reality, banks are defaultable and $\lambda$ is, essentially, a credit spread process (liquidity spreads are typically on the order of a handful of basis points while banks funding spreads can run into the hundreds of basis points).

Collateral means cash or liquid assets that are posted to guarantee a netted set of transactions against defaults of the counterparties. Collateral may include VM, which tracks the mark-to-market of the client portfolio and is typically rehypothecable, and IM set as a cushion against gap risk, i.e., the risk of slippage between the mark-to-market of the portfolio and its VM during the liquidation period that follows a default. Moreover, by regulation, the bank needs to earmark capital at risk $(\mathrm{CR} \geq 0)$ devoted to coping with unexpected losses. For simplicity, we assume cash only collateral and CR. Posted collateral is remunerated at the risk-free rate $r$ by the receiving party.

As explained in section 3.2 in Albanese, Caenazzo, and Crépey (2017), the bank can use reserve capital and CR as VM, but not as IM. As CR can only be used for funding VM and because this feature is our main funding focus in this paper, we assume no IM in the following. IM can be added without harm as done in Albanese et al. (2019, section A) (see also Albanese, Caenazzo, and Crépey (2017)).

Remark 2.2. Bilateral transactions under SIMM, which is the IM standard for noncleared (vanilla) derivatives, are subject to even higher levels of initial margining than centrally cleared transactions. For such transactions, the MVA issue dominates the FVA one. The FVA is dominant in the case of deals with clients subject to little or partial collateralization, so that the bank needs to borrow in order to fund the VM required on its hedge.

We denote by $P$ the difference between the VM posted by the bank on its market hedge and the VM received by the bank on its client portfolio. The market exposure of the bank to the default of its client is modeled as $Q_{t} \boldsymbol{\delta}(d t)$, where $\boldsymbol{\delta}$ is a Dirac measure at the (positive)
default time of the client and where $Q$ is some $\mathbb{F}$ optional nonnegative loss process of the bank given the client default.

Example 2.1. Let MtM (for mark-to-market) denote the value process of the client portfolio ignoring counterparty risk and risky funding costs, i.e., the conditional expectation of the future contractual cash flows $D$ promised by the client to the bank, discounted at the risk-free rate. In line with our assumption of a perfectly collateralized market hedge of its client portfolio by the bank, the bank posts MtM as VM on its hedge. Let VM denote the VM exchanged between the client and the bank, counted positively when received by the bank. Let $R$ denote the recovery rate of the client. Then, assuming instantaneous liquidation of the bank portfolio in case the client defaults (and no IMs), we have ${ }^{4}$

$$
\begin{aligned}
& P=(\mathrm{MtM}-\mathrm{VM}) \text { killed at the default time of the client } \\
& Q=(1-R)\left(\mathrm{MtM}+D-D_{-}-\mathrm{VM}\right)^{+} .
\end{aligned}
$$

The jump ( $D-D_{-}$) of the contractual cash flows contributes to the exposure $Q$ of the bank to the default of its client, consistent with the fact that it fails to be paid by the client if the latter defaults.

Remark 2.3. The XVA setup of this paper can be readily extended to a bank engaged in bilateral trade portfolios with several clients, as considered in Albanese, Caenazzo, and Crépey (2017) and Albanese et al. (2019, section A), by summing the $Q \boldsymbol{\delta}$ and $P$ processes over the different clients of the bank in all equations.

All our XVA processes are a priori sought for as semimartingales (meant componentwise in the multivariate case), which will all happen to be nonnegative special semimartingales.
2.1. Contra-assets valuation. We assume that the process $\lambda P^{+}$is in $\mathcal{H}_{1}$ and that $\int_{0}^{T} Q_{s} \boldsymbol{\delta}(d s)$ is $\mathbb{P}$ integrable. We define

$$
\begin{equation*}
\mathrm{CVA}_{t}=\mathbb{E}_{t}\left[\int_{t}^{T} \beta_{t}^{-1} \beta_{s} Q_{s} \boldsymbol{\delta}(d s)\right], t \in[0, T] \tag{4}
\end{equation*}
$$

which is a nonnegative special semimartingale, like $\int_{0}^{*} \beta_{t}^{-1} \beta_{s} Q_{s} \boldsymbol{\delta}(d s)$ (cf. He, Wang, and Yan (1992, Corollary 11.26)). We define the local martingale

$$
\begin{equation*}
\mu=m(\mathrm{CVA})+m(Q \cdot \boldsymbol{\delta}) . \tag{5}
\end{equation*}
$$

Hence, for $t \in[0, T]$,

$$
\mathrm{CVA}_{0}+\beta \cdot \mu_{t}=m(\beta \mathrm{CVA})+m(\beta Q \cdot \boldsymbol{\delta})=\beta \mathrm{CVA}+\beta Q \cdot \boldsymbol{\delta}=\mathbb{E}_{t}\left[\int_{0}^{T} \beta_{s} Q_{s} \boldsymbol{\delta}(d s)\right],
$$

by (4). That is, $\mu_{0}=\mathrm{CVA}_{0}$ (recall that the client default time is positive, hence, $Q \cdot \boldsymbol{\delta}_{0}=0$ ) and, for $t \in[0, T]$,

$$
\beta_{t} d \mu_{t}=d(\beta \mathrm{CVA})_{t}+\beta_{t} Q_{t} \boldsymbol{\delta}(d t),
$$

i.e.,

$$
\begin{equation*}
d \mu_{t}=d \mathrm{CVA}_{t}-r_{t} \mathrm{CVA}_{t} d t+Q_{t} \boldsymbol{\delta}(d t) \tag{6}
\end{equation*}
$$

[^4]Remark 2.4. If $\int_{0}^{T} Q_{s} \boldsymbol{\delta}(d s)$ is $\mathbb{P}$ square integrable, then CVA is in $\mathcal{S}_{2}$, by (4), and so is $\mu$, by (6) (having assumed $r$ bounded).

The nonnegative contra-asset value process $(\mathrm{CA} \geq 0)$ corresponds to the conditionally expected future counterparty default and risky funding losses, i.e., $\mathrm{CA}=\mathrm{CVA}+\mathrm{FVA}$, computed under the risk-neutral pricing measure $\mathbb{P}$. A well-grounded, refined specification of the FVA is in fact part of the objectives of the paper. The ensuing CA amount is assumed to be charged to the client at time 0 by the CVA desk and the FVA desk (or Treasury) of the bank, which put it into a reserve capital (RC) account dedicated to cope with these expected losses as time goes on. From an accounting perspective, the CVA and the FVA represent special liabilities, which arise from the feedback impact of counterparty risk on financial receivables to the bank, hence the name of contra-assets for the cash flows valued by CA (cf. Figure 1 in Albanese et al. (2019, section 2)).

We assume that all the losses and earnings of the bank are marked to the model and realize immediately. In particular, the RC amount is reset to its theoretical CA value at all times. RC can also be used as VM. Accounting for this feature and since $\mathrm{RC}=\mathrm{CA}$, the VM borrowing needs of the bank are reduced from $P^{+}$otherwise to $(P-\mathrm{CA})^{+}$.

Remark 2.5. The identity $\mathrm{RC}=\mathrm{CA}$ does of course not mean that RC and contra-assets are one and the same thing: RC is the amount on a cash account which is part of the assets of the bank (and can, in particular, be pledged as VM), whereas the contra-assets valuation CA is the matching liability in the bank balance sheet (cf. again Figure 1 in Albanese et al. (2019, section 2)).

On top of $\mathrm{RC}, \mathrm{CR}(\mathrm{CR} \geq 0)$ can also be used by the bank for its funding purposes (provided the bank pays to its shareholders a risk-free rate on CR that they would gain by depositing it otherwise). The funding needs of the bank are then reduced further from $(P-\mathrm{CA})^{+}$to $(P-\mathrm{CA}-\mathrm{CR})^{+}$.

Rephrasing the above qualitative descriptions in mathematical terms, the trading loss process $L$ of the bank satisfies

$$
\begin{align*}
& L_{0}=z \text { and, for } t \in(0, T, \\
& d L_{t}=\underbrace{Q_{t} \boldsymbol{\delta}(d t)}_{\text {loss in case of client default }} \\
& +\underbrace{\left(r_{t}+\lambda_{t}\right)\left(P_{t}-\mathrm{CA}_{t}-\mathrm{CR}_{t}\right)^{+} d t}_{\text {portfolio funding costs }}-\underbrace{r_{t}\left(P_{t}-\mathrm{CA}_{t}-\mathrm{CR}_{t}\right)^{-} d t}_{\text {portfolio funding benefits }} \\
& +\underbrace{r_{t} \mathrm{CR}_{t} d t}_{\text {remuneration to shareholders for the CR funding resource }}  \tag{7}\\
& -\underbrace{r_{t} P_{t} d t}_{\text {remuneration of the collateral between the bank and the client }} \\
& +\underbrace{d \mathrm{CA}_{t}}_{\text {appreciation of the contra-assets of the bank }} \\
& =d \mathrm{CA}_{t}-r_{t} \mathrm{CA}_{t} d t+Q_{t} \boldsymbol{\delta}(d t)+\lambda_{t}\left(P_{t}-\mathrm{CA}_{t}-\mathrm{CR}_{t}\right)^{+} d t
\end{align*}
$$

or, equivalently (recall $\beta=e^{-\int_{0} r_{t} d t}$ ),

$$
\begin{equation*}
\beta_{t} d L_{t}=d\left(\beta_{t} \mathrm{CA}_{t}\right)+\beta_{t} Q_{t} \boldsymbol{\delta}(d t)+\beta_{t} \lambda_{t}\left(P_{t}-\mathrm{CA}_{t}-\mathrm{CR}_{t}\right)^{+} d t \tag{8}
\end{equation*}
$$

(and $L$ is constant beyond $T$ ). Note that the above dynamics for $L$ are well-defined under the postulated integrability conditions on $P$ and $Q$, as well as $\mathrm{CA} \geq 0, \mathrm{CR} \geq 0$. In the case where the stochastic integral $\int_{0}^{0} \beta_{t} d L_{t}$ is a uniformly integrable martingale, formula (8), together with our integrability conditions on $P$ and $Q$ and a terminal condition $\mathrm{CA}_{T}=0$, lead to the following equation for the CA process: For $t \in[0, T]$,

$$
\begin{equation*}
\mathrm{CA}_{t}=\underbrace{\mathbb{E}_{t}\left[\int_{t}^{T} \beta_{t}^{-1} \beta_{s} Q_{s} \boldsymbol{\delta}(d s)\right]}_{\mathrm{CVA}_{t}}+\underbrace{\mathbb{E}_{t}\left[\int_{t}^{T} \beta_{t}^{-1} \beta_{s} \lambda_{s}\left(P_{s}-\mathrm{CA}_{s}-\mathrm{CR}_{s}\right)^{+} d s\right]}_{\mathrm{FVA}_{t}} \tag{9}
\end{equation*}
$$

2.2. KVA. As visible in (9), CR is FVA reducing. However, under a cost-of-capital XVA approach, before being FVA reducing as a side effect, CR entails a charge for the bank, which is the cost, called KVA, of remunerating shareholders at some hurdle rate (dividend rate) for their CR. This cost is sourced from the client at time 0 and put into the so-called risk margin (RM) account, from where it is gradually released, as dividends, to bank shareholders.

Under our continuous reset assumption on all bank accounts, the RM amount is reset to its theoretical KVA value at all times. Moreover, under our XVA approach, the RM is loss absorbing (see sections 3.6 and 6.2 in Albanese and Crépey (2019)), so that it is part of CR. Hence, the inequality

$$
\begin{equation*}
\mathrm{KVA} \leq \mathrm{CR} \tag{10}
\end{equation*}
$$

holds and shareholder CR only corresponds to the difference (CR - KVA). Assuming a constant hurdle rate $h \geq 0$ (including a risk-free remuneration of the RM account to shareholders), this leads to the KVA equation

$$
-d \mathrm{KVA}_{t}+r_{t} \mathrm{KVA}_{t} d t=h\left(\mathrm{CR}_{t}-\mathrm{KVA}_{t}\right) d t-d \nu_{t}, \quad 0 \leq t \leq T,
$$

for some local martingale $\nu$ or, equivalently, ${ }^{5}$

$$
-d\left(\beta_{t} \mathrm{KVA}_{t}\right)=h \beta_{t}\left(\mathrm{CR}_{t}-\mathrm{KVA}_{t}\right) d t-\beta_{t} d \nu_{t}, \quad 0 \leq t \leq T
$$

In the case where the stochastic integral $\int_{0}^{0} \beta_{t} d \nu_{t}$ is a uniformly integrable martingale, the above formula, together with the terminal condition $\mathrm{KVA}_{T}=0$ (as the portfolio expires at time $T$ ), leads to

$$
\begin{equation*}
0 \leq \beta_{t} \mathrm{KVA}_{t}=h \mathbb{E}_{t}\left[\int_{t}^{T} \beta_{s}\left(\mathrm{CR}_{s}-\mathrm{KVA}_{s}\right) d s\right], \quad 0 \leq t \leq T \tag{11}
\end{equation*}
$$

[^5]2.3. Economic capital and CR specifications. Economic capital is the (minimum) level of CR that a regulator would like to see on an economic, structural basis. We use the following dynamic specification of economical capital:
\[

$$
\begin{equation*}
\mathbb{E S}_{t}\left(\int_{t}^{t+1} \beta_{t}^{-1} \beta_{s} d L_{s}\right), t \in[0, T] \tag{12}
\end{equation*}
$$

\]

where $L$ is the trading loss and profit process of the bank in (8) and where $\mathbb{E S}_{t}(\ell)$ denotes the $\mathfrak{F}_{t}$ conditional expected shortfall, at some level $\alpha$ (e.g., $\alpha=97.5 \%$ ), of an $\mathfrak{F}_{T}$ measurable, $\mathbb{P}$ integrable random variable $\ell$. That is, denoting by $q_{t}^{a}(\ell)$ the $\left(\mathfrak{F}_{t}, \mathbb{P}\right)$ conditional value at risk (left quantile) of level $a$ of $\ell$ (cf. Artzner et al. (1999)): For $t \leq T$,

$$
\begin{align*}
\mathbb{E}_{t}(\ell) & =(1-\alpha)^{-1} \int_{\alpha}^{1} q_{t}^{a}(\ell) d a \\
& =\inf _{x \in \mathbb{R}}\left(x+(1-\alpha)^{-1} \mathbb{E}_{t}\left[(\ell-x)^{+}\right]\right)  \tag{13}\\
& =\sup \left\{\mathbb{E}_{t}[\ell \chi] ; \chi \text { is } \mathfrak{F}_{T} \text { measurable, } 0 \leq \chi \leq(1-\alpha)^{-1}, \text { and } \mathbb{E}_{t}[\chi]=1\right\}
\end{align*}
$$

Note that we will only deal with martingale loss and profit processes $L$ and therefore centered loss variables $\ell$ fer which $\mathbb{E S} \mathbb{S}_{t}(\ell) \geq 0$ holds in view of the third line in (13). Moreover, for any $\mathfrak{F}_{T}$ measurable, $\mathbb{P}$ integrable random variables $\ell$ and $\ell^{\prime}$, we have (cf. Barrera et al. (2019, Lemma 6.10, eq. (6.20)) and its proof)

$$
\begin{equation*}
\left|\mathbb{E S}_{t}(\ell)-\mathbb{E S}_{t}\left(\ell^{\prime}\right)\right| \leq(1-\alpha)^{-1} \mathbb{E}_{t}\left[\left|\ell-\ell^{\prime}\right|\right], 0 \leq t \leq T \tag{14}
\end{equation*}
$$

We will need to plug economic capital processes such as (12) into the coefficient of BSDEs. Before doing so (see, e.g., (27)), let us show that any such economic capital process is the image, in the sense of (16) below, of a predictable process. Namely, for every (raw, nonnecessarily adapted) process $\Lambda$ admitting ${ }^{6}$ an $(\mathbb{F}, \mathbb{P})$ predictable projection ${ }^{p} \Lambda$, let

$$
\begin{equation*}
\mathrm{EC}(\Lambda)=\inf _{\text {rationals } k}\left(k+(1-\alpha)^{-1 p}\left[(\Lambda-k)^{+}\right]\right) \tag{15}
\end{equation*}
$$

Lemma 2.1. For every raw process $\Lambda$ endowed with an $(\mathbb{F}, \mathbb{P})$ predictable projection, we have, for every constant $t \geq 0$,

$$
\begin{equation*}
\operatorname{EC}(\Lambda)_{t}=\mathbb{E S}_{t}\left(\Lambda_{t}\right) \tag{16}
\end{equation*}
$$

Proof. By definition,

$$
\operatorname{EC}(\Lambda)_{t}=\inf _{\text {rationals } k}\left(k+(1-\alpha)^{-1} \mathbb{E}_{t-}\left[\left(\Lambda_{t}-k\right)^{+}\right]\right)
$$

By Jacod (1979, Lemma 4.48),

$$
\mathfrak{F}_{t}=\mathfrak{F}_{t-} \vee \sigma\left(\Delta_{t} X: \text { all uniformly integrable martingale } X\right) .
$$

[^6]By the martingale representation property (3), the martingales can have no jump at the predictable stopping time $t$. Consequently, $\mathfrak{F}_{t}=\mathfrak{F}_{t-}$ and

$$
\operatorname{EC}(\Lambda)_{t}=\inf _{\text {rationals } k}\left(k+(1-\alpha)^{-1} \mathbb{E}_{t}\left[\left(\Lambda_{t}-k\right)^{+}\right]\right)=\mathbb{E S}_{t}\left(\Lambda_{t}\right)
$$

by the second line in (13).
Hence, by application of Lemma 2.1 (to $\Lambda=\int_{.+1}^{-1} \beta^{-1} \beta_{s} d L_{s}$ in (16)), it is harmless to use processes such as $\mathbb{E S}$. $\left(\int^{+1} \beta^{-1} \beta_{s} d L_{s}\right)$ (cf. (12)) instead of the corresponding predictable expression $\operatorname{EC}\left(\int_{.}^{+1} \beta^{-1} \beta_{s} d L_{s}\right)$ in our XVA equations below (where $L$ is a local martingale, as implied by the requirements for a solution to these XVA equations).

In view of inequality (10) and of the risk admissibility condition

$$
\mathrm{CR} \geq \mathbb{E S} .\left(\int_{.}^{+1} \beta_{\cdot}^{-1} \beta_{s} d L_{s}\right)
$$

(see before (12)), we adopt the following "minimal" specification of CR:

$$
\begin{equation*}
\mathrm{CR}=\max \left(\mathbb{E S}\left(\int_{.}^{+1} \beta^{-1} \beta_{s} d L_{s}\right), \mathrm{KVA}\right) \tag{17}
\end{equation*}
$$

Table 1 reviews the main protagonists of the cost-of-capital XVA approach.
Table 1
Main financial acronyms and place where they are introduced conceptually and/or specified mathematically in the paper, as relevant.

| CA | Contra-assets valuation | Section 2.1 and (27) |
| :--- | :--- | :--- |
| CR | Capital at risk | Paragraph following Remark 2.1 and (17) |
| CVA | Credit valuation adjustment | First paragraph of section 1 and (4) |
| EC | Economic capital | (12) and preceding lines |
| FVA | Funding valuation adjustment | First paragraph of section 1 and (9), (27) |
| KVA | Capital valuation adjustment | Section 2.2 and (11), (27) |
| MtM | Mark-to-market | Example 2.1 |
| RC | Reserve capital | Section 2.1, Remark 2.5 in particular |
| RM | Risk margin | Section 2.2 |
| XVA | Generic "X" valuation adjustment | First paragraph of section 1 |

3. A general ABSDE well-posedness result. In this section, we extend to ABSDEs the (square integrable part of the) monotone coefficient BSDE results in Kruse and Popier (2016).

Let there be given a map $\rho$ from $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ into the space of $\mathbb{F}$ predictable processes satisfying the following.

Assumption 3.1. There exists a constant $c_{\rho}$ such that, for any $(Y, Z, U),\left(Y^{\prime}, Z^{\prime}, U^{\prime}\right)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$, for any $t \in[0, T]$,

$$
\begin{align*}
& \left|\rho_{t}(Y, Z, U)-\rho_{t}\left(Y^{\prime}, Z^{\prime}, U^{\prime}\right)\right|^{2} \\
& \quad \leq c_{\rho}^{2} \mathbb{E}_{t}\left[\sup _{t \leq s \leq T}\left|Y_{s}-Y_{s}^{\prime}\right|^{2}+\int_{t}^{T}\left(\left|Z_{s}-Z_{s}^{\prime}\right|^{2}+\left|U_{s}-U_{s}^{\prime}\right|_{s}^{2}\right) d s\right] . \tag{18}
\end{align*}
$$

Let there additionally be given an $\mathfrak{F}_{T}$ measurable terminal condition $\xi \in \mathbb{R}^{l}$ and a $\mathfrak{P} \otimes$ $\mathfrak{B}\left(\mathbb{R}^{l}\right) \otimes \mathfrak{B}\left(\mathbb{R}^{l \otimes d}\right) \otimes \mathfrak{B}\left(\mathcal{L}_{0}^{l}\right) \otimes \mathfrak{B}(\mathbb{R})$ measurable coefficient $f$, such that we have the following.

Assumption 3.2.
(i) The function $y \mapsto f(t, y, z, u, \varrho)$ is continuous. Moreover, there exists a positive constant $c_{m}$ such that

$$
\left\langle f(t, y, z, u, \varrho)-f\left(t, y^{\prime}, z, u, \varrho\right), y-y^{\prime}\right\rangle \leq c_{m}\left|y-y^{\prime}\right|^{2}
$$

(ii) There exists a positive constant $c_{f}$ such that

$$
\left|f(t, y, z, u, \varrho)-f\left(t, y, z^{\prime}, u^{\prime}, \varrho^{\prime}\right)\right| \leq c_{f}\left(\left|z-z^{\prime}\right|+\left|u-u^{\prime}\right|_{t}+\left|\varrho-\varrho^{\prime}\right|\right)
$$

(iii) The processes $\sup _{|y| \leq c}|f(\cdot, y, 0,0, \rho \cdot(0,0,0))-f(\cdot, 0,0,0, \rho \cdot(0,0,0))|($ for every $c>0)$, as well as $|f(\cdot, 0,0,0, \rho \cdot(0,0,0))|^{2}$, are in $\mathcal{H}_{1}$.
(iv) $\mathbb{E}\left[|\xi|^{2}\right]<+\infty$.

We consider the following $l$-variate ABSDE with data $\xi, f, \rho$ :

$$
\left\{\begin{align*}
Y_{T} & =\xi \text { and, for } t \leq T  \tag{19}\\
-d Y_{t} & =f\left(t, Y_{t}, Z_{t}, U_{t}, \rho_{t}(Y, Z, U)\right) d t-Z_{t} d B_{t}-\int_{E} U_{t}(e) M(d t, d e)
\end{align*}\right.
$$

to be solved for $(Y, Z, U)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$.
Lemma 3.1. If $(Y, Z, U)$ is a solution in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ to the $\operatorname{ABSDE}$ (19), then

$$
\begin{equation*}
\|Y\|_{\mathcal{S}_{2}^{l}}^{2}+\|Z\|_{\mathcal{H}_{2}^{l}}^{2}+\|U\|_{\widehat{\mathcal{H}}_{2}^{l}}^{2} \leq C \mathbb{E}\left[|\xi|^{2}+\int_{0}^{T}|f(s, 0,0,0, \rho \cdot(0,0,0))|^{2} d s\right] \tag{20}
\end{equation*}
$$

Proof. See section A.1.
Lemma 3.2. For any given $(X, V, W)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$, the BSDE

$$
\left\{\begin{array}{l}
Y_{T}=\xi \text { and, for } t \leq T  \tag{21}\\
-d Y_{t}=f\left(t, Y_{t}, Z_{t}, U_{t}, \rho_{t}(X, V, W)\right) d t-Z_{t} d B_{t}-\int_{E} U_{t}(e) M(d t, d e)
\end{array}\right.
$$

has a unique solution $(Y, Z, U)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$.
Proof. Given $(X, V, W)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l},(21)$ for $(Y, Z, U)$ is a monotone BSDE, which is well-posed in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$, by Kruse and Popier (2016, Theorem 1).

Theorem 3.1. The $A B S D E$ (19) has a unique solution $(Y, Z, U)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$, which is the limit in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$, with a geometrical convergence rate, of the Picard iteration defined by $\left(Y^{(0)}, Z^{(0)}, U^{(0)}\right)=(0,0,0)$ and, for $n \geq 1$,

$$
\left\{\begin{array}{l}
Y_{T}^{(n)}=\xi \text { and, for } t \leq T  \tag{22}\\
-d Y_{t}^{(n)}=f\left(t, Y_{t}^{(n)}, Z_{t}^{(n)}, U_{t}^{(n)}, \rho_{t}\left(Y^{(n-1)}, Z^{(n-1)}, U^{(n-1)}\right)\right) d t \\
\\
\quad-Z_{t}^{(n)} d B_{t}-\int_{E} U_{t}^{(n)}(e) M(d t, d e)
\end{array}\right.
$$

Proof. Lemma 3.2 allows defining a map $\Phi$ from $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ into itself by $\Phi(X, V, W)=$ $(Y, Z, U)$. Classical BSDE computations detailed in section A. 2 show that $\Phi$ is a contraction on the space $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ (endowed with a suitably weighted norm still making it a Banach space), so that $\Phi$ has a unique fix-point $(Y, Z, U)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$. Hence, (19) has a solution $(Y, Z, U)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$, unique by Lemma 3.1.

Moreover, the sequence ( $Y^{(n)}, Z^{(n)}, U^{(n)}$ ), is well-defined, by iterated application of Lemma 3.2, and the majoration (69) implies the stated geometrical convergence of this sequence to $(Y, Z, U)$.

Remark 3.1. In the case where $f$ is not only monotone but Lipschitz in $y$, then one can check likewise that the explicit Picard iteration, with $Y_{t}^{(n-1)}$ instead of $Y_{t}^{(n)}$ as second argument of $f$ in (22), also converges in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ to $(Y, Z, U)$.
3.1. A special case. In this subsection, we assume that $\rho(Y, Z, U)$ only depends on $(Z, U)$ through $\int_{0}^{s} Z_{s} d B_{s}+\int_{0} \int_{E} U_{s}(e) M(d s, d e)$ and that $f$ does not depend on $(z, u)$. More precisely,

$$
\rho(Y, Z, U)=\bar{\rho}\left(Y, Y_{0}+\int_{0} Z_{s} d B_{s}+\int_{0} \int_{E} U_{s}(e) M(d s, d e)\right), f(t, y, z, u, \varrho)=\bar{f}(t, y, \varrho),
$$

where $\bar{f}$ satisfies the amended form of Assumption 3.2 obtained by replacing all missing arguments by 0 there, while $\bar{\rho}$ is a map from $\mathcal{S}_{2}^{l} \times \mathcal{M}_{2}^{l}$ into the space of $\mathbb{F}$ predictable processes, which satisfies the following.

Assumption 3.3. There exists a constant $c_{\bar{\rho}}$ such that, for any $(Y, N),\left(Y^{\prime}, N^{\prime}\right)$ in $\mathcal{S}_{2}^{l} \times \mathcal{M}_{2}^{l}$, for any $t \in[0, T]$,

$$
\begin{align*}
& \left|\bar{\rho}_{t}(Y, N)-\bar{\rho}_{t}\left(Y^{\prime}, N^{\prime}\right)\right|^{2} \\
& \quad \leq c_{\bar{\rho}}^{2} \mathbb{E}_{t}\left[\sup _{t \leq s \leq T}\left|Y_{s}-Y_{s}^{\prime}\right|^{2}+\sum_{k=1}^{l}\left(\left\langle N^{k}-\left(N^{\prime}\right)^{k}\right\rangle_{T}-\left\langle N^{k}-\left(N^{\prime}\right)^{k}\right\rangle_{t}\right)\right] . \tag{23}
\end{align*}
$$

Then $\left|\rho_{t}(Y, Z, U)-\rho_{t}\left(Y^{\prime}, Z^{\prime}, U^{\prime}\right)\right|^{2}$ in (18) is rewritten as $\mid \bar{\rho}(Y, N)-\bar{\rho}\left(Y^{\prime},\left.N^{\prime}\right|^{2}\right.$, where

$$
\begin{aligned}
& N=Y_{0}+\int_{0} Z_{s} d B_{s}+\int_{0} \int_{E} U_{s}(e) M(d s, d e) \text { and } \\
& N^{\prime}=Y_{0}^{\prime}+\int_{0} Z_{s}^{\prime} d B_{s}+\int_{0} \int_{E} U_{s}^{\prime}(e) M(d s, d e)
\end{aligned}
$$

are such that

$$
\sum_{k=1}^{l}\left(\left\langle N^{k}-\left(N^{\prime}\right)^{k}\right\rangle_{T}-\left\langle N^{k}-\left(N^{\prime}\right)^{k}\right\rangle_{t}\right)=\int_{t}^{T}\left(\left|Z_{s}-Z_{s}^{\prime}\right|^{2}+\left|U_{s}-U_{s}^{\prime}\right|_{s}^{2}\right) d s
$$

Hence (23) is rewritten as (18) (for $c_{\rho}$ there equal to $c_{\bar{\rho}}$ ), which shows that Assumption 3.1 is satisfied. Moreover, by inspection,

1. Assumption 3.2 holds;
2. the $\operatorname{ABSDE}$ (19) for $(Y, Z, U)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ is equivalent, via the martingale representation property (3), to the following equation to be solved for a (special) semimartingale $Y$ in $\mathcal{S}_{2}^{l}$ with $m(Y)$ in $\mathcal{S}_{2}^{l}$ :

$$
\begin{equation*}
Y_{t}=\mathbb{E}_{t}\left[\xi+\int_{t}^{T} \bar{f}\left(s, Y_{s}, \bar{\rho}_{s}(Y, m(Y))\right) d s\right], t \leq T \tag{24}
\end{equation*}
$$

which is in turn equivalent to the following system of equations for a (special) semimartingale $Y$ in $\mathcal{S}_{2}^{l}$ and a martingale $N(=m(Y))$ in $\mathcal{S}_{2}^{l}$ :

$$
\begin{align*}
& N_{0}=Y_{0} \text { and, for } t \in(0, T] \\
& d N_{t}=d Y_{t}-\bar{f}\left(t, Y_{t}, \bar{\rho}_{t}(Y, N)\right) d t \\
& Y_{t}=\mathbb{E}_{t}\left[\xi+\int_{t}^{T} \bar{f}\left(s, Y_{s}, \bar{\rho}_{s}(Y, N)\right) d s\right] \tag{25}
\end{align*}
$$

3. the Picard iteration of Theorem 3.1 for $(Y, Z, U)$ is equivalent to the following Picard iteration for $(Y, N=m(Y))$ in $(25): Y^{(0)}=N^{(0)}=0$ and, for $n \geq 1$,

$$
\begin{align*}
& Y_{t}^{(n)}=\mathbb{E}_{t}\left[\xi+\int_{t}^{T} \bar{f}\left(s, Y_{s}^{(n)}, \bar{\rho}_{s}\left(Y^{(n-1)}, N^{(n-1)}\right)\right) d s\right], 0 \leq t \leq T \\
& N_{0}^{(n)}=Y_{0}^{(n)} \text { and, for } t \in(0, T]  \tag{26}\\
& d N_{t}^{(n)}=d Y_{t}^{(n)}-\bar{f}\left(t, Y_{t}^{(n)}, \bar{\rho}_{t}\left(Y^{(n-1)}, N^{(n-1)}\right)\right) d t
\end{align*}
$$

4. The XVA equations in the case of a default-free bank are well-posed.
4.1. The (FVA,KVA) ABSDE system. Note that $L$ only intervenes via CR as per (17) in the FVA and KVA equations (7), (9), (11), and (17). As a consequence, using in this paragraph a superscript $z$ in reference to the value of the initial condition $z$ for $L$ in (7), if $\left(L^{z}, \mathrm{FVA}^{z}, \mathrm{KVA}^{z}\right)$ solves the FVA and KVA equations, then, for every real $u,\left(L^{z}+u-\right.$ $z, \mathrm{FVA}^{z}, \mathrm{KVA}^{z}$ ) solves the same equations with $z$ replaced by $u$ in (7). Hence, the value of the initial condition $z$ for $L$ is immaterial in the XVA problem.

Hereafter, in order to be in line with the convention that the martingale part of a special semimartingale $S$ starts from $S_{0}$ (in this case we target $L=m(\mathrm{CA}+Q \cdot \boldsymbol{\delta})$ ), we set $z=\mathrm{CA}_{0}$. The FVA and KVA problem (7), (9), (11), and (17) is then rewritten as (with CA used as a shorthand for CVA + FVA, hence $L=\mu+m(\mathrm{FVA})$, by (5))

$$
\begin{align*}
L_{0}= & \mathrm{CA}_{0} \text { and, for } t \in(0, T] \\
= & d \mathrm{CA}_{t}-r_{t} \mathrm{CA}_{t} d t+Q_{t} \boldsymbol{\delta}(d t) \\
& +\lambda_{t}\left(P_{t}-\mathrm{CA}_{t}-\max \left(\mathbb{E S}_{t}\left(\int_{t}^{t+1} \beta_{t}^{-1} \beta_{s} d L_{s}\right), \mathrm{KVA}_{t}\right)\right)^{+} d t \\
\mathrm{FVA}_{t}= & \mathbb{E}_{t} \int_{t}^{T} \beta_{t}^{-1} \beta_{s} \lambda_{s}\left(P_{s}-\mathrm{CA}_{s}-\max \left(\mathbb{E}_{s}\left(\int_{s}^{s+1} \beta_{s}^{-1} \beta_{u} d L_{u}\right), \mathrm{KVA}_{s}\right)\right)^{+} d s  \tag{27}\\
\mathrm{KVA}_{t}= & h \mathbb{E}_{t}\left[\int_{t}^{T} \beta_{t}^{-1} \beta_{s}\left(\mathbb{E S}_{s}\left(\int_{s}^{s+1} \beta_{s}^{-1} \beta_{u} d L_{u}\right)-\mathrm{KVA}_{s}\right)^{+} d s\right]
\end{align*}
$$

As will be detailed in the proof of Theorem 4.1 below, the system (27) belongs formally to the class of equations considered in section 3.1, in the form (25), for $Y=(\mathrm{FVA}, \mathrm{KVA})$ and its martingale part $N$ (having noted that $L=\mu+m$ (FVA) in (27)).

Remark 4.1. As $\lambda$ is bounded (and $h$ is constant), the system (27) is Lipschitz in ( $\beta$ FVA, $\beta \mathrm{KVA}$ ), irrespective of the boundedness assumption made on $r$. However, we prefer to look at the system in terms of (FVA, KVA), even if this is at the cost of assuming $r$ bounded (which is not a real restriction in practice, as explained at the beginning of section 2). Indeed, in Markov setups with numerical solutions in mind, working with ( $\beta \mathrm{FVA}, \beta \mathrm{KVA}$ ) would require us to introduce an additional factor process (at least mathematically) to account for the path dependence induced by $\beta$.

Definition 4.1. We denote as a square integrable solution to (27), any special semimartingale solution (componentwise) (FVA, KVA) (with $L=\mu+m(\mathrm{FVA})$ in (27)), with components (i.e., the XVA processes themselves) and their martingale parts in $\mathcal{S}_{2}$.

### 4.2. The (FVA,KVA) ABSDE system (27) is well-posed.

Lemma 4.1. The map defined by, for $(Y, N)$ in $\mathcal{S}_{2}^{2} \times \mathcal{M}_{2}^{2}$,

$$
\begin{equation*}
\bar{\rho}_{t}(Y, N)=\mathbb{E S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s}\left(d \mu_{s}+d N_{s}^{1}\right)\right], t \in[0, T] \tag{28}
\end{equation*}
$$

(which only depends on $(Y, N)$ through the first component $N^{1}$ of $N$ ), satisfies Assumption 3.3 (for $l=2$ there).

Proof. By (14), we have for any $(Y, N),\left(Y^{\prime}, N^{\prime}\right)$ in $\mathcal{S}_{2}^{2} \times \mathcal{M}_{2}^{2}$,

$$
\begin{aligned}
\left|\bar{\rho}_{t}(Y, N)-\bar{\rho}_{t}\left(Y^{\prime}, N^{\prime}\right)\right|^{2} & \leq(1-\alpha)^{-2}\left(\mathbb{E}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d\left(N_{s}^{1}-\left(N^{\prime}\right)_{s}^{1}\right)\right]\right)^{2} \\
& \leq(1-\alpha)^{-2} \mathbb{E}_{t}\left[\left(\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d\left(N_{s}^{1}-\left(N^{\prime}\right)_{s}^{1}\right)\right)^{2}\right]
\end{aligned}
$$

by the (conditional) Jensen inequality. Moreover, as a local martingale in $\mathcal{S}_{2}$, the process $\left(N^{1}-\left(N^{\prime}\right)^{1}\right)$ is a square integrable martingale. As $\beta_{t}^{-1} \beta_{s}, t \leq s \leq T$, is bounded, the process $s \mapsto \int_{t}^{\prime} \beta_{t}^{-1} \beta_{s} d\left(N_{s}^{1}-\left(N^{\prime}\right)_{s}^{1}\right)$ is therefore a square integrable martingale over $[t, T]$. The (conditional) Burkholder inequality applied to this process then yields

$$
\left|\bar{\rho}_{t}(Y, N)-\bar{\rho}_{t}\left(Y^{\prime}, N^{\prime}\right)\right|^{2} \leq(1-\alpha)^{-2} C \mathbb{E}_{t}\left(\left\langle N^{1}-\left(N^{\prime}\right)^{1}\right\rangle_{T}-\left\langle N^{1}-\left(N^{\prime}\right)^{1}\right\rangle_{t}\right)
$$

In particular, Assumption 3.3 holds (with $l=2$ there).
We introduce the following Picard iteration, where CVA and $\mu$ are given by (4) and (5): $\mathrm{FVA}^{(0)}=\mathrm{KVA}^{(0)}=0, L^{(0)}=\mu$ and, for $n \geq 1,\left(\mathrm{FVA}^{(n)}, \mathrm{KVA}^{(n)}\right)$ given as the unique square
integrable solution ${ }^{7}$ to, with $\mathrm{CA}^{(n)}$ as a shorthand for $\mathrm{CVA}+\mathrm{FVA}^{(n)}$,

$$
\begin{aligned}
\mathrm{KVA}_{t}^{(n)} & =h \mathbb{E}_{t} \int_{t}^{T} \beta_{t}^{-1} \beta_{s}\left(\mathbb{E S}_{s}\left(\int_{s}^{s+1} \beta_{s}^{-1} \beta_{u} d L_{u}^{(n-1)}\right)-\mathrm{KVA}_{s}^{(n)}\right)^{+} d s, t \in[0, T] \\
\mathrm{FVA}_{t}^{(n)} & =\mathbb{E}_{t} \int_{t}^{T} \beta_{t}^{-1} \beta_{s} \lambda_{s}\left(P_{s}-\mathrm{CA}_{s}^{(n)}\right.
\end{aligned}
$$

$$
\begin{equation*}
\left.-\max \left(\mathbb{E S}_{s}\left(\int_{s}^{s+1} \beta_{s}^{-1} \beta_{u} d L_{u}^{(n-1)}\right), \mathrm{KVA}_{s}^{(n)}\right)\right)^{+} d s, t \in[0, T] \tag{29}
\end{equation*}
$$

$$
\begin{aligned}
L_{0}^{(n)}= & \mathrm{CA}_{0}^{(n)} \text { and, for } t \in[0, T], d L_{t}^{(n)}=Q_{t} \boldsymbol{\delta}(d t)+d \mathrm{CA}_{t}^{(n)}-r_{t} \mathrm{CA}_{t}^{(n)} d t \\
& +\lambda_{t}\left(P_{t}-\mathrm{CA}_{t}^{(n)}-\max \left(\mathbb{E} \mathbb{S}_{t}\left(\int_{t}^{t+1} \beta_{t}^{-1} \beta_{s} d L_{s}^{(n-1)}\right), \mathrm{KVA}_{t}^{(n)}\right)\right)^{+} d t
\end{aligned}
$$

Theorem 4.1. Suppose ( $r$ and $\lambda$ are bounded), $\left(P^{+}\right)^{2} \in \mathcal{H}_{1}$ and $\int_{0}^{T} Q_{s} \boldsymbol{\delta}(d s)$ is $\mathbb{P}$ square integrable. Then the system (27) admits a unique square integrable solution (FVA, KVA), which is also the limit in $\mathcal{S}_{2}^{2}$, with a geometrical convergence rate, of the sequence $\left(\mathrm{FVA}^{(n)}, \mathrm{KVA}^{(n)}\right)_{n \in \mathbb{N}}$ in (29). Moreover, $L=\mu+m(\mathrm{FVA}) \in \mathcal{M}_{2}$.

Proof. The system (27) is nothing but the integral form of (25) with
$l=2, \bar{\rho}$ as per $(28), \xi=(0,0)^{\top}$, and $\bar{f}=\left(\bar{f}_{1}, \bar{f}_{2}\right)^{\top}$ such that, for any $t, y=\left(y_{1}, y_{2}\right)^{\top}, \varrho$,

$$
\begin{align*}
& \bar{f}_{1}=\bar{f}_{1}(t, y, \varrho)=\lambda_{t}\left(P_{t}-\mathrm{CVA}_{t}-y_{1}-\max \left(\varrho, y_{2}\right)\right)^{+}-r_{t} y_{1}  \tag{30}\\
& \bar{f}_{2}=\bar{f}_{2}(t, y, \varrho)=h\left(\varrho-y_{2}\right)^{+}-r_{t} y_{2}
\end{align*}
$$

As shown by Lemma 4.1, $\bar{\rho}$ in (28) satisfies Assumption 3.3. Moreover, for any $t \in[0, T]$, for $y=\left(y_{1}, y_{2}\right)^{\top}$ and $y^{\prime}=\left(y_{1}^{\prime}, y_{2}^{\prime}\right)^{\top}$ in $\mathbb{R}^{2}$, as $r$ and $\lambda$ are bounded from below, we have

$$
\begin{aligned}
& \left\langle\bar{f}(t, y, \varrho)-\bar{f}\left(t, y^{\prime}, \varrho\right), y-y^{\prime}\right\rangle \\
& =\left(\bar{f}_{1}(t, y, \varrho)-\bar{f}_{1}\left(t, y^{\prime}, \varrho\right)\right)\left(y_{1}-y_{1}^{\prime}\right)+\left(\bar{f}_{2}(t, y, \varrho)-\bar{f}_{2}\left(t, y^{\prime}, \varrho\right)\right)\left(y_{2}-y_{2}^{\prime}\right) \\
& =\lambda_{t}\left(\left(P_{t}-\mathrm{CVA}_{t}-y_{1}-\max \left(\varrho, y_{2}\right)\right)^{+}-\left(P_{t}-\mathrm{CVA}_{t}-y_{1}^{\prime}-\max \left(\varrho, y_{2}^{\prime}\right)\right)^{+}\right)\left(y_{1}-y_{1}^{\prime}\right) \\
& \quad \quad-r_{t}\left(y_{1}-y_{1}^{\prime}\right)^{2}+h\left(\left(\varrho-y_{2}\right)^{+}-\left(\varrho-y_{2}^{\prime}\right)^{+}\right)\left(y_{2}-y_{2}^{\prime}\right)-r_{t}\left(y_{2}-y_{2}^{\prime}\right)^{2} \\
& \leq C\left|y-y^{\prime}\right|^{2}
\end{aligned}
$$

Hence, the ABSDE coefficient $\bar{f}$ in (30) satisfies the monotonicity condition of Assumption $3.2(\mathrm{i})$ (it is in fact even Lipschitz, as our processes $r$ and $\lambda$ are actually assumed bounded). Next, for any arguments $t, y=\left(y_{1}, y_{2}\right)^{\top}$, and $\varrho, \varrho^{\prime}$, we have

$$
\begin{aligned}
& \left|\bar{f}(t, y, \varrho)-\bar{f}\left(t, y, \varrho^{\prime}\right)\right|^{2} \\
& =\left(\bar{f}_{1}(t, y, \varrho)-\bar{f}_{1}\left(t, y, \varrho^{\prime}\right)\right)^{2}+\left(\bar{f}_{2}(t, y, \varrho)-\bar{f}_{2}\left(t, y, \varrho^{\prime}\right)\right)^{2} \\
& =\quad \lambda_{t}^{2}\left|\left(P_{t}-\mathrm{CVA}_{t}-y_{1}-\max \left(\varrho, y_{2}\right)\right)^{+}-\left(P_{t}-\mathrm{CVA}_{t}-y_{1}-\max \left(\varrho^{\prime}, y_{2}\right)\right)^{+}\right|^{2} \\
& \quad \quad \quad+h^{2}\left|\left(\varrho-y_{2}\right)^{+}-\left(\varrho^{\prime}-y_{2}\right)^{+}\right|^{2} \leq C\left|\varrho-\varrho^{\prime}\right|^{2}
\end{aligned}
$$

[^7]which implies that Assumption 3.2(ii) holds. Finally, writing
$$
\widehat{\rho}_{t}=\bar{\rho}_{t}\left((0,0)^{\boldsymbol{\top}},(0,0)^{\boldsymbol{\top}}\right)=\mathbb{E S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right]
$$
we have, for each $t$ and $y=\left(y_{1}, y_{2}\right)^{\top}$ (simply denoting by 0 the origin of $\mathbb{R}^{2}$ in what follows)
\[

$$
\begin{align*}
\bar{f}_{1}\left(t, y, \widehat{\rho}_{t}\right)= & \lambda_{t}\left(P_{t}-\mathrm{CVA}_{t}-y_{1}\right. \\
& \left.\quad-\max \left(\mathbb{E} \mathbb{S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right], y_{2}\right)\right)^{+}-r_{t} y_{1}, \\
\bar{f}_{1}\left(t, 0, \widehat{\rho}_{t}\right)= & \lambda_{t}\left(P_{t}-\mathrm{CVA}_{t}-\max \left(\mathbb{E S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right]\right)\right)^{+}  \tag{31}\\
\bar{f}_{2}\left(t, y, \widehat{\rho}_{t}\right)= & h\left(\mathbb{E S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right]-y_{2}\right)^{+}-r_{t} y_{2} \\
\bar{f}_{2}\left(t, 0, \widehat{\rho}_{t}\right)= & h \mathbb{E S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right]
\end{align*}
$$
\]

hence,

$$
\left|\bar{f}\left(t, y, \widehat{\rho}_{t}\right)-\bar{f}\left(t, 0, \widehat{\rho}_{t}\right)\right| \leq\left(\left|\lambda_{t}\right|+\left|r_{t}\right|\right)\left(\left|y_{1}\right|+\left|y_{2}\right|\right)+h\left|y_{2}\right| .
$$

By assumption, $r$ and $\lambda$ are bounded and $\left(P^{+}\right)^{2} \in \mathcal{H}_{1}$. In addition,

$$
\begin{aligned}
\left\|\hat{\rho}^{2}\right\|_{\mathcal{H}_{1}} & =\mathbb{E}\left[\int_{0}^{T}\left(\mathbb{E} \mathbb{S}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right]\right)^{2} d t\right] \\
& \leq(1-\alpha)^{-2} \mathbb{E}\left[\int_{0}^{T}\left(\mathbb{E}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \mu_{s}\right]\right)^{2} d t\right] \\
& \leq(1-\alpha)^{-2} \mathbb{E}\left[\int_{0}^{T}\left(\mathbb{E}_{t}\left[\int_{t}^{(t+1) \wedge T} \beta_{t}^{-2} \beta_{s}^{2} d[\mu]_{s}\right]\right) d t\right] .
\end{aligned}
$$

Hence, since $\beta_{t}^{-2} \beta_{s}^{2} \leq e^{C T}$,

$$
\left\|\hat{\rho}^{2}\right\|_{\mathcal{H}_{1}} \leq(1-\alpha)^{-2} e^{C T} \mathbb{E}\left[\int_{0}^{T} \int_{0}^{T} d[\mu]_{s} d t\right] \leq(1-\alpha)^{-2} e^{C T} T \mathbb{E}\left[\mu_{T}^{2}\right]
$$

which is finite, by Remark 2.4 (having assumed that $\int_{0}^{T} Q_{s} \boldsymbol{\delta}(d s)$ is $\mathbb{P}$ square integrable). From the above (starting with (31) onward), we deduce that Assumption 3.2(iii) is satisfied. Hence, an application of Theorem 3.1, via the observations 1 and 2 in section 3.1, shows that (25) with data (30) has a unique solution in $\mathcal{S}_{2}^{2} \times \mathcal{M}_{2}^{2}$, i.e., the (FVA, KVA) system (27) admits a unique square integrable solution.

The well-posedness of the Picard iteration (29) among square integrable solutions $\left(\mathrm{FVA}^{(n)}, \mathrm{KVA}^{(n)}\right)$, for each $n$, and the geometrical convergence of $\left(\mathrm{FVA}^{(n)}, \mathrm{KVA}^{(n)}\right)$, when $n \rightarrow \infty$, follow through the observation 3 in section 3.1, from the second part of Theorem 3.1.

Finally, we have $L=\mu+m(\mathrm{FVA})$ in (27), where $\mu$ is in $\mathcal{M}_{2}$, by Remark 2.4, and so is $m(\mathrm{FVA})$, by Definition 4.1 of a square integrable solution to (27).

Remark 4.2. As $L=\mu+m(\mathrm{FVA}) \in \mathcal{S}_{2}$ (last statement in Theorem 4.1) and $\beta$ is bounded, the stochastic integral $\int_{0}^{u} \beta_{t} d L_{t}$ is a uniformly integrable martingale, in line with the corresponding requirement before (9).
5. Invariance times transfer properties. What precedes was done in the theoretical case of a default-free bank. However, in reality, banks are defaultable and counterparty risk is related to cash flows or valuations linked to either counterparty default or the default of the bank itself. In particular, the bank funding spread $\lambda$, which we introduced in the above as liquidity, is essentially related to the credit spread of the bank (see Remark 2.1). Hence it is crucial to understand the case of a defaultable bank.

Then, as detailed in Albanese and Crépey (2019, section 4.1) and recalled in what follows, we consider not only one pricing basis, i.e., $(\mathbb{F}, \mathbb{P})$ above, but actually two, i.e., also $(\mathbb{G}, \mathbb{Q})$ below. These pricing bases are connected by suitable consistency conditions (B) and (A), meaning that the bank default time $\tau$ is an invariance time in the sense of Crépey and Song (2017) (see Appendix B for a survey of the main results there).

In this transition section, which is general and of independent interest, we establish a transfer between expectations and martingale properties in the original and changed stochastic bases, assuming an invariance time $\tau$ endowed with an intensity and a positive Azéma supermartingale.
5.1. Reduction of filtration setup. We suppose that a complete and right-continuous filtration $\mathbb{G}=\left(\mathfrak{G}_{t}\right)_{t \in \mathbb{R}_{+}}$of sub- $\sigma$ fields of $\mathfrak{A}$ is an enlargement of our previous filtration $\mathbb{F}=\left(\mathfrak{F}_{t}\right)_{t \in \mathbb{R}_{+}}$, such that we have the following.

Condition (B). ${ }^{8} \forall t \geq 0$ and $B \in \mathfrak{G}_{t}, \exists B^{\prime} \in \mathfrak{F}_{t}$ such that $B \cap\{t<\tau\}=B^{\prime} \cap\{t<\tau\}$.
This holds, in particular, under a standard progressive enlargement of the filtration setup, i.e., when $\mathfrak{G}_{t}=\mathfrak{F}_{t} \vee \sigma(\tau \wedge t), t \geq 0$.

As seen in Crépey and Song (2017, Lemma 2.2), under the condition (B), any $\mathbb{G}$ optional process $Y$ admits an $\mathbb{F}$ optional reduction $Y^{\prime}$ such that $\mathbb{1}_{[0, \tau)} Y=\mathbb{1}_{[0, \tau)} Y^{\prime}$; any $\mathbb{G}$ predictable process $Y$ admits an $\mathbb{F}$ predictable reduction $Y^{\prime}$ such that $\mathbb{1}_{(0, \tau]} Y=\mathbb{1}_{(0, \tau]} Y^{\prime}$.

For any process $Y$, we denote by $Y^{\tau}=\mathbb{1}_{[0, \tau)} Y+\mathbb{1}_{[\tau,+\infty)} Y_{\tau}$ and ${ }^{9}$ by $Y^{\tau-}=\mathbb{1}_{[0, \tau)} Y+$ $\mathbb{1}_{[\tau,+\infty)} Y_{\tau-}$ the processes $Y$ stopped at and before $\tau$, respectively. A process $Y$ is said to be stopped at $\tau$ if $Y=Y^{\tau}$ and stopped before $\tau$ if $Y=Y^{\tau-}$.

Given a positive constant $T$, we assume further that a full model probability measure $\mathbb{Q}$ equivalent to $\mathbb{P}$ on $\mathfrak{F}_{T}$ satisfies the following.

Condition (A). For any $(\mathbb{F}, \mathbb{P})$ local martingale $P$ on $[0, T], P^{\tau-}$ is a $(\mathbb{G}, \mathbb{Q})$ local martingale on $[0, T]$.

The full model probability measure $\mathbb{Q}$ is then interpreted as the prevailing risk-neutral pricing measure, whereas the reduced stochastic pricing basis $(\mathbb{F}, \mathbb{P})$ plays a technical role analogous to the one of the survival probability measure associated with $\mathbb{Q}$ (cf. Schönbucher (2004)), while avoiding the singularity issue of the latter (see Crépey and Song (2017, section 4.2)).

[^8]The Azéma supermartingale $S$ of $\tau$ is defined as $S_{t}=\mathbb{Q}\left(\tau>t \mid \mathfrak{F}_{t}\right), t>0$.
Remark 5.1. The situation where $\mathbb{P}=\mathbb{Q}$ in condition (A) corresponds to the special case where $S$ is nonincreasing and predictable. The subcase where $S$ is also continuous corresponds to the class of pseudostopping times with the avoidance property (see Nikeghbali and Yor (2005) and Crépey and Song (2017, section 4.1)). Pseudostopping times include Cox times, the family of default times most commonly used in the credit and counterparty risk literatures; see, e.g., Bielecki, Jeanblanc, and Rutkowski (2009, Chapter 3), Brigo and Pallavicini (2014), or Bichuch, Capponi, and Sturm (2018).

We assume $\mathrm{S}_{T}>0$ a.s., so that, by Crépey and Song (2017, Lemma 2.3),

$$
\begin{equation*}
\text { Two } \mathbb{F} \text { optional processes that coincide before } \tau \text { coincide on }[0, T] \text {. } \tag{32}
\end{equation*}
$$

In particular, $\mathbb{F}$ optional reductions are uniquely defined on $[0, T]$.
Finally, we assume that $\tau$ has a $(\mathbb{G}, \mathbb{Q})$ intensity $\gamma$, i.e., the $(\mathbb{G}, \mathbb{Q})$ compensator of $\mathbb{1}_{[\tau, \infty)}$ is given as $\int_{0}^{\cdot} \gamma_{t} d t$ for some $\mathbb{G}$ predictable process $\gamma$ (vanishing beyond time $\tau$ ). Summarizing, we suppose in the remainder of sections 5-6 the following.

Condition (C). The conditions $(\mathrm{B})$ and $(\mathrm{A})$ are satisfied, $\mathrm{S}_{T}>0$, and $\tau$ has a $(\mathbb{G}, \mathbb{Q})$ intensity $\gamma$.

We then have the following converse to the condition (A).
Lemma 5.1. For any $(\mathbb{G}, \mathbb{Q})$ local martingale $M$ stopped before $\tau$ (i.e., such that $M=$ $\left.M^{\tau-}\right), M^{\prime}$ is an $(\mathbb{F}, \mathbb{P})$ local martingale on $[0, T]$.

Proof. The process $\mathrm{S}_{-} \cdot M^{\prime}+\left[\mathrm{S}, M^{\prime}\right]$ is an $(\mathbb{F}, \mathbb{Q})$ local martingale on $\left\{\mathrm{S}_{-}>0\right\}$, by (71), where $\left\{S_{-}>0\right\}=\{\mathrm{S}>0\} \supseteq[0, T]$ by (75) and our assumption that $\mathrm{S}_{T}>0$. The conclusion then follows from (76).
5.2. Expectation transfer formulas. The following result provides an extension of the classical credit risk intensity pricing formulas (see, e.g., Bielecki, Jeanblanc, and Rutkowski (2009, Chapter 3)) to the invariance time setup, i.e., beyond the basic immersion setup, where $(\mathbb{F}, \mathbb{P}=\mathbb{Q})$ local martingales are $(\mathbb{G}, \mathbb{Q})$ local martingales without jump at $\tau$.

We denote the $\mathbb{Q}$ expectation by $\mathbb{E}$, whereas the $\mathbb{P}$ expectation is denoted as before by $\mathbb{E}$.
Theorem 5.1. For any $\mathbb{F}$ stopping time $\sigma \leq T$ and $\mathfrak{F}_{\sigma}$ measurable nonnegative random variable $\chi$, for any $\mathbb{F}$ predictable nonnegative process $K$, for any $\mathbb{F}$ optional nondecreasing process A starting from 0, we have, respectively,

$$
\begin{align*}
& \widetilde{\mathbb{E}}\left[\chi \mathbb{1}_{\{\sigma<\tau\}}\right]=\mathbb{E}\left[\chi e^{-\int_{0}^{\sigma} \gamma_{s}^{\prime} d s}\right]  \tag{33}\\
& \widetilde{\mathbb{E}}\left[K_{\tau} \mathbb{1}_{\{\tau \leq T\}}\right]=\mathbb{E}\left[\int_{0}^{T} K_{s} e^{-\int_{0}^{s} \gamma_{u}^{\prime} d u} \gamma_{s}^{\prime} d s\right]  \tag{34}\\
& \widetilde{\mathbb{E}}\left[A_{T}^{\tau-}\right]=\mathbb{E}\left[\int_{0}^{T} e^{-\int_{0}^{s} \gamma_{u}^{\prime} d u} d A_{s}\right] \tag{35}
\end{align*}
$$

Proof. Let $\Gamma=\int_{0}^{r} \gamma_{s}^{\prime} d s$. Since $\chi$ is $\mathfrak{F}_{\sigma}$ measurable, $\widetilde{\mathbb{E}}\left[\chi \mathbb{1}_{\{\sigma<\tau\}}\right]=\widetilde{\mathbb{E}}\left[\chi \mathrm{S}_{\sigma}\right]$. As $S_{T}>0$, (75) implies that $\left\{S_{-}>0\right\}=\left\{{ }^{p} S>0\right\} \supseteq[0, T]$. Then (72) implies that

$$
\begin{equation*}
\widetilde{\mathbb{E}}\left[\chi S_{\sigma}\right]=\widetilde{\mathbb{E}}\left[\chi \mathrm{S}_{0} \mathcal{E}\left(-\frac{1}{\mathrm{~S}_{-}} \cdot \mathrm{D}\right)_{\sigma} \mathcal{E}\left(\frac{1}{p \mathrm{~S}} \cdot \mathrm{Q}\right)_{\sigma}\right]=\mathbb{E}\left[\chi \mathrm{S}_{0} \mathcal{E}\left(-\frac{1}{\mathrm{~S}_{-}} \cdot \mathrm{D}\right)_{\sigma}\right] \tag{36}
\end{equation*}
$$

by (74). In view of (73), we obtain (33).
For (34), we compute

$$
\begin{aligned}
\widetilde{\mathbb{E}}\left[K_{\tau} \mathbb{1}_{\{\tau \leq T\}}\right] & =\widetilde{\mathbb{E}}\left[\int_{0}^{T} K_{s} \mathbb{1}_{\{s \leq \tau\}} \gamma_{s}^{\prime} d s\right]=\int_{0}^{T} \widetilde{\mathbb{E}}\left[K_{s} \mathbb{1}_{\{s<\tau\}} \gamma_{s}^{\prime}\right] d s \\
& =\int_{0}^{T} \mathbb{E}\left[K_{s} e^{-\Gamma_{s}} \gamma_{s}^{\prime}\right] d s=\mathbb{E}\left[\int_{0}^{T} K_{s} e^{-\Gamma_{s}} \gamma_{s}^{\prime} d s\right],
\end{aligned}
$$

where (33) was used for passing to the second line.
Regarding (35), an application of (34) yields

$$
\widetilde{\mathbb{E}}\left[A_{\tau-} \mathbb{1}_{\{\tau \leq T\}}\right]=\mathbb{E}\left[\int_{0}^{T} A_{s} e^{-\Gamma_{s}} \gamma_{s}^{\prime} d s\right]=-\mathbb{E}\left[A_{T} e^{-\Gamma_{T}}\right]+\mathbb{E}\left[\int_{0}^{T} e^{-\Gamma_{s}} d A_{s}\right]
$$

Using (33), we deduce

$$
\begin{aligned}
\widetilde{\mathbb{E}}\left[A_{T}^{\tau-}\right] & =\widetilde{\mathbb{E}}\left[A_{T} \mathbb{1}_{\{T<\tau\}}\right]+\widetilde{\mathbb{E}}\left[A_{\tau-} \mathbb{1}_{\{\tau \leq T\}}\right] \\
& =\mathbb{E}\left[A_{T} e^{-\Gamma_{T}}\right]-\mathbb{E}\left[A_{T} e^{\left.-\Gamma_{T}\right]+\mathbb{E}}\left[\int_{0}^{T} e^{-\Gamma_{s}} d A_{s}\right]\right. \\
& =\mathbb{E}\left[\int_{0}^{T} e^{-\Gamma_{s}} d A_{s}\right] .
\end{aligned}
$$

Denoting $X_{t}^{*}=\sup _{t \in[0, T]}\left|X_{s}\right|$, we introduce the space $\widetilde{\mathcal{S}}_{2}^{l}\left(\widetilde{Y}^{(s i m p l y}\right.$ renoted $\widetilde{\mathcal{S}}_{2}$ when $l=1$ ) of the $\mathbb{R}^{l}$ valued càdlàg $\mathbb{G}$ adapted processes $\tilde{Y}$ such that $\widetilde{Y}=\widetilde{Y}^{\tau-}$ and

$$
\begin{equation*}
\|\widetilde{Y}\|_{\widetilde{\mathcal{S}}_{2}^{l}}^{2}=\widetilde{\mathbb{E}}\left[\left|\widetilde{Y}_{0}\right|^{2}+\int_{0}^{T} \mathbb{1}_{\{s<\tau\}} e^{\int_{0}^{s} \gamma_{u} d u} d\left(\widetilde{Y}^{*}\right)_{s}^{2}\right]<\infty . \tag{37}
\end{equation*}
$$

Lemma 5.2. For any real valued càdlàg $\mathbb{F}$ adapted process $Y$, we have

$$
\begin{equation*}
\|Y\|_{\mathcal{S}_{2}}=\|Y\|_{\tilde{\mathcal{S}}_{2}} \tag{38}
\end{equation*}
$$

Proof. For $A=\int_{0} e^{\int_{0}^{s} \gamma_{u}^{\prime} d u} d\left(Y^{*}\right)_{s}^{2}$, the expectation transfer formula (35) yields

$$
\widetilde{\mathbb{E}}\left[\int_{0}^{T} \mathbb{1}_{\{s<\tau\}} e^{\int_{0}^{s} \gamma_{u} d u} d\left(Y^{*}\right)_{s}^{2}\right]=\mathbb{E}\left[\left(Y^{*}\right)_{T}^{2}\right]-\mathbb{E}\left[Y_{0}^{2}\right]=\mathbb{E}\left[\left(Y^{*}\right)_{T}^{2}\right]-\widetilde{\mathbb{E}}\left[Y_{0}^{2}\right]
$$

(noting that $\mathbb{P}$ and $\mathbb{Q}$ coincide on $\mathfrak{F}_{0}$, by (74)), which is (38).
6. The realistic case of a defaultable bank. We now revisit the analysis of section 2 in the realistic case of a defaultable bank, with default time $\tau$ satisfying the condition (C) of section 5.1. We can then identify $\lambda$, the funding spread of the bank, with the instantaneous credit spread process $(1-\bar{R}) \gamma^{\prime}$ of the bank, where $\bar{R}$ is a constant recovery rate (cf. Remark 2.1). The time horizon of the XVA problem is now $\bar{\tau}=\tau \wedge T$, where $T$ is the final maturity of the derivative portfolio of the bank.

Accounting for the defaultability of the bank, a key distinction appears between the cash flows received by the bank prior $\tau$ and the cash flows received by the bank during the default resolution period starting at $\tau$. Indeed, the first stream of cash flows affects the bank shareholders, whereas the second stream of cash flows only affects bondholders. For accepting a new deal, bank shareholders need to be at least indifferent given the first stream of cash flows only. Accordingly, as can be seen from the general CVA, FVA, and KVA equations (10), (11), and (18)-(20) in Albanese and Crépey (2019), ${ }^{10}$ the "recipe" to obtain the XVA equations of a defaultable bank under $(\mathbb{G}, \mathbb{Q})$ is to stop all cash flows before $\tau$ in the $(\mathbb{F}, \mathbb{P})$ equations (or, more precisely but to the same result, to stop all cash flows before $\tau$ in the extension of these equations that also include the cash flows received by the bank from time $\tau$ onward). See below for concrete illustrations.

In all $(\mathbb{G}, \mathbb{Q})$ equations, we write $\widetilde{\text { XVA }}$ for the corresponding (defaultable bank) XVA process stopped before $\tau$. Likewise, we denote by $\widetilde{L}$ (instead of $L$ before) the now $\mathbb{G}$ adapted bank trading loss process $L$ stopped before $\tau$.
6.1. CVA. The differential form of the CVA equation (4) is written as

$$
\begin{align*}
& \mathrm{CVA}_{T}=0 \text { and, for } t \in(0, T], \\
& d \mathrm{CVA}_{t}=r_{t} \mathrm{CVA}_{t} d t-Q_{t} \boldsymbol{\delta}(d t)+d \mu_{t}  \tag{39}\\
& \text { for some }(\mathbb{F}, \mathbb{P}) \text { martingale } \mu \text { in } \mathcal{S}_{2} .
\end{align*}
$$

On the other hand, in the present defaultable bank setup, the differential form of the general CVA equation (10) in Albanese and Crépey (2019) is

$$
\begin{align*}
& \widetilde{\mathrm{CVA}}_{T}=0 \text { on }\{T<\tau\} \text { and, for } t \in[0, \bar{\tau}], \\
& d \widetilde{\mathrm{CVA}}_{t}=r_{t} \widetilde{\mathrm{CVA}}_{t} d t-\mathbb{1}_{\{t<\tau\}} Q_{t} \boldsymbol{\delta}(d t)+d \widetilde{\mu}_{t} \tag{40}
\end{align*}
$$

for some $(\mathbb{G}, \mathbb{Q})$ martingale $\widetilde{\mu}$ in $\widetilde{\mathcal{S}}_{2}$.
Observe that (40) is nothing but the $(\mathbb{G}, \mathbb{Q})$ equation obtained by formally ignoring all cash flows from $\tau$ onward in (39).

Now, (39) and (40) are in fact equivalent via $\mathbb{F}$ (optional) reduction ${ }^{\prime}$ (cf. section 5.1).
Proposition 6.1. Equation (39) for CVA in $\mathcal{S}_{2}$ and (40) for $\widetilde{\mathrm{CVA}}$ in $\widetilde{\mathcal{S}}_{2}$ are equivalent via $\mathbb{F}$ optional reduction, i.e., through the relation $\mathrm{CVA}=\widetilde{\mathrm{CVA}}{ }^{\prime}$.

If $\int_{0}^{T} Q_{s} \boldsymbol{\delta}(d s)$ is $\mathbb{P}$ square integrable, then (39) for CVA in $\mathcal{S}_{2}$ has the unique solution (4).

[^9]Proof. Assuming that $(\widetilde{\mathrm{CVA}}, \widetilde{\mu}) \in \widetilde{\mathcal{S}}_{2}^{2}$ satisfies (40), then $(\mathrm{CVA}, \mu)=\left(\widetilde{\mathrm{CVA}}^{\prime}, \widetilde{\mu}^{\prime}\right) \in \mathcal{S}_{2}^{2}$, by (38), and (CVA $\mu$ ) thus defined satisfies the second line in (39) on $[0, \bar{\tau}]$, hence, on $[0, T]$, by (32), while the martingale condition in the third line holds by Lemma 5.1. Moreover, taking the $\left(\mathfrak{F}_{T}, \mathbb{Q}\right)$ conditional expectation of the first line in (40) yields

$$
0=\widetilde{\mathbb{E}}\left[\widetilde{\mathrm{CVA}}_{T} \mathbb{1}_{\{T<\tau\}} \mid \mathfrak{F}_{T}\right]=\widetilde{\mathbb{E}}\left[\widetilde{\mathrm{CVA}}_{T}^{\prime} \mathbb{1}_{\{T<\tau\}} \mid \mathfrak{F}_{T}\right]=\widetilde{\mathrm{CVA}}_{T}^{\prime} \mathrm{S}_{T},
$$

where $\mathrm{S}_{T}>0$, hence, $\mathrm{CVA}_{T}=\widetilde{\mathrm{CVA}}_{T}^{\prime}=0$.
Conversely, assuming that (CVA, $\mu$ ) $\in \mathcal{S}_{2}^{2}$ satisfies (39), then

$$
(\widetilde{\mathrm{CVA}}, \widetilde{\mu})=\left(\mathrm{CVA}^{\tau-}, \mu^{\tau-}\right) \in \widetilde{\mathcal{S}}_{2}^{2},
$$

by (38), and $(\widetilde{\mathrm{CVA}}, \widetilde{\mu})=\left(\mathrm{CVA}^{\tau-}, \mu^{\tau-}\right)$ satisfies the conditions of the first two lines in (40), while $\widetilde{\mu}=\mu^{\tau-}$ is a $(\mathbb{G}, \mathbb{Q})$ local martingale, by virtue of the condition (A).

This proves the first part in Proposition 6.1. The second part is an immediate consequence of Remark 2.4.
6.2. FVA and KVA. The differential form of the (FVA, KVA) ABSDE system (27) is written as (with CA as a shorthand for CVA + FVA):

$$
\begin{align*}
\mathrm{CA}_{T}= & 0 \text { and, for } t \in(0, T] \\
d \mathrm{CA}_{t} & =r_{t} \mathrm{CA}_{t} d t-Q_{t} \boldsymbol{\delta}(d t) \\
& -\lambda_{t}\left(P_{t}-\max \left(\mathbb{E S}_{t}\left(\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d L_{s}\right), \mathrm{KVA}_{t}\right)-\mathrm{CA}_{t}\right)^{+} d t+d L_{t} \tag{41}
\end{align*}
$$

for some $(\mathbb{F}, \mathbb{P})$ martingale $L$ in $\mathcal{S}_{2}$
( $L$ was interpreted in section 2 as the trading loss of a default-free bank), along with

$$
\mathrm{KVA}_{T}=0 \text { and, for } t \in(0, T],
$$

$$
\begin{equation*}
d \mathrm{KVA}_{t}=\left(r_{t}+h\right) \mathrm{KVA}_{t} d t-h \max \left(\mathbb{E S}_{t}\left(\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d L_{s}\right), \mathrm{KVA}_{t}\right) d t+d \nu_{t} \tag{42}
\end{equation*}
$$

for some $(\mathbb{F}, \mathbb{P})$ martingale $\nu$ in $\mathcal{S}_{2}$.
On the other hand, in the present defaultable bank setup, the general $\widetilde{\text { FVA }}$ and $\widetilde{\text { KVA }}$ equations (11) and (18)-(20) in Albanese and Crépey (2019) yield the following system of $\widetilde{\mathrm{XVA}}$ equations (with $\widetilde{\mathrm{CA}}$ used as a shorthand for $\widetilde{\mathrm{CVA}}+\widetilde{\mathrm{FVA}}$ ):

$$
\begin{align*}
{\widetilde{\mathrm{CA}_{T}}}= & 0 \text { on }\{T<\tau\} \text { and, for } t \in(0, \bar{\tau}] \\
d \widetilde{\mathrm{CA}}_{t} & =r_{t} \widetilde{\mathrm{CA}}_{t} d t-\mathbb{1}_{\{t<\tau\}} Q_{t} \boldsymbol{\delta}(d t) \\
& -\lambda_{t}\left(P_{t}-\widetilde{\mathrm{CA}_{t}}-\max \left(\mathbb{E S}_{t}\left(\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \widetilde{L}_{s}^{\prime}\right), \widetilde{\mathrm{KVA}}_{t}\right)\right)^{+} d t+d \widetilde{L}_{t} \tag{43}
\end{align*}
$$

for some $(\mathbb{G}, \mathbb{Q})$ martingale $\widetilde{L}$ in $\widetilde{\mathcal{S}}_{2}$
( $\widetilde{L}$ is interpreted as the bank shareholders' trading loss of the now defaultable bank), along with

$$
\begin{align*}
& \widetilde{\mathrm{KVA}}_{T}=0 \text { on }\{T<\tau\} \text { and, for } t \in(0, \bar{\tau}] \\
& d \widetilde{\mathrm{KVA}}_{t}=\left(r_{t}+h\right) \widetilde{\mathrm{KVA}}_{t} d t-h \max \left(\mathbb{E S}_{t}\left(\int_{t}^{(t+1) \wedge T} \beta_{t}^{-1} \beta_{s} d \widetilde{L}_{s}^{\prime}\right), \widetilde{\mathrm{KVA}}_{t}\right) d t+d \widetilde{\nu}_{t}  \tag{44}\\
& \text { for some }(\mathbb{G}, \mathbb{Q}) \text { martingale } \widetilde{\nu} \text { in } \widetilde{\mathcal{S}}_{2} .
\end{align*}
$$

Remark 6.1. We recall from Albanese and Crépey (2019, section 4.3) that capital calculations are typically performed "on a going-concern basis," i.e., disregarding the default of the bank itself. In view of the comments preceding Remark 5.1, this grounds our specification $\mathbb{E} \mathbb{S}_{t}\left(\int_{t}^{t+1} \beta_{t}^{-1} \beta_{s} d \widetilde{L}_{s}^{\prime}\right)$ in (43)-(44) for the economic capital of a defaultable bank.

Lemma 6.1. Equations (41)-(42) for (FVA, KVA) in $\mathcal{S}_{2}^{2}$ and (43)-(44) for $(\widetilde{\mathrm{FVA}}, \widetilde{\mathrm{KVA}})$ in $\widetilde{\mathcal{S}}_{2}^{2}$ are equivalent via $\mathbb{F}$ optional reduction, i.e., through the relation $(\mathrm{FVA}, \mathrm{KVA})=$ ( $\left.\widetilde{\mathrm{FVA}^{\prime}}, \widetilde{\mathrm{KVA}}^{\prime}\right)$.

Proof. The proof is similar to the one of the first part of Proposition 6.1, hence omitted.
Square integrable solutions to the $(\widetilde{\mathrm{FVA}}, \widetilde{\mathrm{KVA}})$ equations are defined in reference to the $\widetilde{\mathcal{S}}_{2}$ space.

Theorem 6.1. Under the assumptions of Theorem 4.1, the $(\widetilde{\mathrm{FVA}}, \widetilde{\mathrm{KVA}})$ and (FVA, KVA) equations are well-posed (and equivalent via $\mathbb{F}$ optional reduction) in their respective spaces of square integrable solutions.

Proof. This is the conclusion of Lemma 6.1 and Theorem 4.1 (first part), also noting that, under the assumptions of Theorem 4.1, the differential and the integral formulations of the (FVA, KVA) system are equivalent in their space of square integrable solutions.

Remark 6.2. Under the assumptions of Theorem 4.1, one also has the equivalence, within square integrable solutions, between the differential formulation (43)-(44) of the ( $\widetilde{(\mathrm{FVA}}, \widetilde{\mathrm{KVA}})$ equations and the corresponding integral formulation (not displayed in the paper for length's sake). This can be shown by the same argument as the one used for concluding the proof of the first part of Theorem 5.1 in Albanese and Crépey (2019).
7. Synthesis. To conclude, we reread the paper upside down, which reveals the overall logic of our XVA approach.

By application of the general CVA, FVA, and KVA equations in Albanese and Crépey (2019, sections 3-4), we derive in section 6 the differential formulations (40) and (43)-(44) of these equations in the specific setup that we want to address in this paper, with possible use of CR as VM by the bank. In particular, the FVA and KVA equations (43)-(44) form a bivariate system of ABSDEs (BSDEs of the McKean type), stopped before the bank default time $\tau$.

Using an enlargement of filtration methodology of independent interest, finalized in section 5 following up on Crépey and Song (2017), we show in section 6 the equivalence between these ABSDEs, natively stated with respect to the "full" pricing model including the bank default
time as a stopping time, and reduced ABSDEs (41)-(42) stated with respect to a smaller filtration and an equivalently changed probability measure.

These reduced equations are none but the ones that were derived in section 2 in the irrealistic case of a default-free bank, at the cost of abusively interpreting the bank funding spread there as liquidity (cf. Remark 2.1). By ignoring the default of the bank, the XVA setup of section 2 was far-fetched and, in a sense, internally inconsistent financially. But the related mathematics of sections $3-4$, showing that the corresponding bank default-free XVA equations were well-posed, are of course valid. Hence, we were able to conclude in Theorem 6.1 that the bank default-prone FVA and KVA equations (43)-(44) are well-posed mathematically.

Putting together Theorems 6.1 and 4.1, we obtain schematically, under the assumptions of Theorem 4.1,

$$
\begin{equation*}
(\widetilde{\mathrm{FVA}}, \widetilde{\mathrm{KVA}}) \sim(\mathrm{FVA}, \mathrm{KVA}) \leftarrow\left(\mathrm{FVA}^{(n)}, \mathrm{KVA}^{(n)}\right) . \tag{45}
\end{equation*}
$$

Moreover, as demonstrated numerically in Albanese, Caenazzo, and Crépey (2017, section 5) and Albanese, Crépey, Hoskinson, and Saadeddine (2019, section 5), the Picard iteration $\left(\mathrm{FVA}^{(n)}, \mathrm{KVA}^{(n)}\right)$ in (29) is amenable to Monte Carlo approximation, including at the scale of a real banking portfolio with hundreds of counterparties and hundreds of thousands of contracts. Hence, the XVA approach of this paper is not only theoretically well-posed, but also workable in practice for a bank.

A more extensive numerical study of the XVA ABSDEs is left for future research.
Appendix A. Details of the ABSDE proofs. In this section, we give the details of the computations used to establish the well-posedness of the $l$-variate ABSDE (19).

First we note that, if ( $Y, Z, U$ ) solves (19), then

$$
\left(\mathcal{Y}_{t}, \mathcal{Z}_{t}, \mathcal{U}_{t}\right)=\left(e^{c_{m} t} Y_{t}, e^{c_{m} t} Z_{t}, e^{c_{m} t} U_{t}\right)
$$

solves, for $0 \leq t \leq T$,

$$
\left\{\begin{align*}
\mathcal{Y}_{t}=e^{c_{m} T} & \xi+\int_{t}^{T}\left(e^{c_{m} s} f\left(s, e^{-c_{m} s} \mathcal{Y}_{s}, e^{-c_{m} s} \mathcal{Z}_{s}, e^{-c_{m} s} \mathcal{U}_{s}, \rho_{s}\left(e^{-c_{m}} \cdot \mathcal{Y}_{.}, e^{-c_{m} \cdot \mathcal{Z} .}, e^{-c_{m}} \cdot \mathcal{U} .\right)\right)\right.  \tag{46}\\
& \left.-c_{m} \mathcal{Y}_{s}\right) d s-\int_{t}^{T} \int_{E} \mathcal{U}_{s}(e) M(d s, d e)-\int_{t}^{T} \mathcal{Z}_{s} d B_{s} .
\end{align*}\right.
$$

Namely, $(\mathcal{Y}, \mathcal{Z}, \mathcal{U})$ solves an $l$-variate $\operatorname{ABSDE}$ of the general form (19), but for modified data that satisfy Assumptions 3.1 and 3.2 with $c_{m}=0$ in the latter. Moreover,

$$
(Y, Z, U) \in \mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l} \Longleftrightarrow(\mathcal{Y}, \mathcal{Z}, \mathcal{U}) \in \mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}
$$

Hence we may and do suppose that $c_{m}=0$ in Assumption 3.2, without loss of generality in what follows.
A.1. Proof of Lemma 3.1. For $t \in[0, T]$, set $\Gamma_{t}=e^{\kappa t}$, for some to-be-determined positive constant $\kappa$. An application of the Itô formula to $\Gamma|Y|^{2}$ yields

$$
\begin{align*}
\Gamma_{t}\left|Y_{t}\right|^{2} & +\int_{t}^{T} \Gamma_{s}\left|Z_{s}\right|^{2} d s+\int_{t}^{T} \Gamma_{s}\left|U_{s}\right|_{s}^{2} d s=\Gamma_{T}\left|Y_{T}\right|^{2}-\kappa \int_{t}^{T} \Gamma_{s}\left|Y_{s}\right|^{2} d s \\
& +2 \int_{t}^{T} \Gamma_{s}\left\langle Y_{s}, f\left(s, Y_{s}, Z_{s}, U_{s}, \rho_{s}(Y, Z, U)\right)\right\rangle d s  \tag{47}\\
& -2 \int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}-\int_{t}^{T} \int_{E} \Gamma_{s}\left(\left|Y_{s-}+U_{s}(e)\right|^{2}-\left|Y_{s-}\right|^{2}\right) M(d s, d e)
\end{align*}
$$

The Burkholder's inequality yields

$$
\begin{gathered}
\mathbb{E}\left[\sup _{t \in[0, T]}\left|\int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}\right|\right] \leq C \mathbb{E}\left[\left(\int_{0}^{T} \Gamma_{s}^{2}\left|Y_{s}\right|^{2}\left|Z_{s}\right|^{2} d s\right) g^{1 / 2}\right] \\
\leq \mathbb{E}\left[\left(\sup _{s \in[0, T]} \Gamma_{s}\left|Y_{s}\right|^{2}\right)^{1 / 2}\left(C \int_{0}^{T} \Gamma_{s}\left|Z_{s}\right|^{2} d s\right)^{1 / 2}\right]
\end{gathered}
$$

Then, by Young's inequality ( $a b \leq \frac{1}{2 \epsilon} a^{2}+\frac{\epsilon}{2} b^{2}$ ) with $\epsilon=4$, we have

$$
\begin{equation*}
\mathbb{E}\left[\sup _{t \in[0, T]}\left|\int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}\right|\right] \leq \frac{1}{8} \mathbb{E}\left(\sup _{s \in[0, T]} \Gamma_{s}\left|Y_{s}\right|^{2}\right)+2 C^{2} \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left|Z_{s}\right|^{2} d s\right] . \tag{48}
\end{equation*}
$$

Similarly, we obtain

$$
\begin{align*}
& \mathbb{E}\left[\sup _{t \in[0, T]}\left|\int_{t}^{T} \Gamma_{s} \int_{E} Y_{s-} U_{s}(e) M(d s, d e)\right|\right] \\
& \leq C \mathbb{E}\left[\left(\int_{0}^{T} \Gamma_{s}^{2}\left|Y_{s-}\right|^{2} \int_{E}\left|U_{s}(e)\right|^{2} j(d s, d e)\right)^{1 / 2}\right]  \tag{49}\\
& \leq \mathbb{E}\left[\left(\sup _{s \in[0, T]} \Gamma_{s}\left|Y_{s}\right|^{2}\right)^{1 / 2}\left(C \int_{0}^{T} \Gamma_{s} \int_{E}\left|U_{s}(e)\right|^{2} j(d s, d e)\right)^{1 / 2}\right] \\
& \leq \frac{1}{8} \mathbb{E}\left(\sup _{s \in[0, T]} \Gamma_{s}\left|Y_{s}\right|^{2}\right)+2 C^{2} \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left|U_{s}\right|_{s}^{2} d s\right]
\end{align*}
$$

In particular, the $\left(\mathfrak{F}_{t}, \mathbb{P}\right)$ conditional expectation of the last line in (47) is equal to 0 .
Moreover, by Assumption 3.2 on $f$ and Young's inequality, we have, for any $\epsilon>0$,

$$
\begin{aligned}
& 2\left\langle Y_{s}, f\left(s, Y_{s}, Z_{s}, U_{s}, \rho_{s}(Y, Z, U)\right)\right\rangle \\
& \leq 2\left|Y_{s}\right|\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|+2 c_{f}\left|Y_{s}\right|\left(\left|Z_{s}\right|+\left|U_{s}\right|_{s}+\left|\rho_{s}(Y, Z, U)-\rho_{s}(0,0,0)\right|\right) \\
& \leq\left(1+8 c_{f}^{2}+c_{f}^{2} \epsilon^{-1}\right)\left|Y_{s}\right|^{2}+\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} \\
& \quad+\frac{1}{4}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right)+\epsilon\left|\rho_{s}(Y, Z, U)-\rho_{s}(0,0,0)\right|^{2} .
\end{aligned}
$$

Besides,

$$
\begin{align*}
& \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\sup _{s \leq u \leq T}\left|Y_{u}\right|^{2}+\int_{s}^{T}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u\right) d s\right] \\
& \leq \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\sup _{s \leq u \leq T}\left|Y_{u}\right|^{2}\right) d s\right]+\mathbb{E}\left[\int_{0}^{T} \int_{s}^{T} \Gamma_{s}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u d s\right] \\
& \leq \mathbb{E}\left[\int_{0}^{T}\left(\sup _{s \leq u \leq T} \Gamma_{u}\left|Y_{u}\right|^{2}\right) d s\right]+\mathbb{E}\left[\int_{0}^{T} \int_{0}^{u} \Gamma_{s}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d s d u\right]  \tag{51}\\
& \leq T \mathbb{E}\left[\sup _{0 \leq u \leq T} \Gamma_{u}\left|Y_{u}\right|^{2}\right]+T \mathbb{E}\left[\int_{0}^{T} \Gamma_{u}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u\right] \\
& \leq T \mathbb{E}\left[\sup _{0 \leq u \leq T} \Gamma_{u}\left|Y_{u}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right]
\end{align*}
$$

In addition, Assumption 3.1 implies

$$
\begin{align*}
& \int_{t}^{T} \Gamma_{s}\left|\rho_{s}(Y, Z, U)-\rho_{s}(0,0,0)\right|^{2} d s \\
& \qquad \leq c_{\rho}^{2} \int_{t}^{T} \Gamma_{s} \mathbb{E}_{s}\left(\sup _{s \leq u \leq T}\left|Y_{u}\right|^{2}+\int_{s}^{T}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u\right) d s \tag{52}
\end{align*}
$$

hence,

$$
\begin{align*}
& \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left|\rho_{s}(Y, Z, U)-\rho_{s}(0,0,0)\right|^{2} d s\right]  \tag{53}\\
& \quad \leq c_{\rho}^{2} T \mathbb{E}\left[\sup _{0 \leq u \leq T} \Gamma_{u}\left|Y_{u}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right]
\end{align*}
$$

Plugging (50) into (47), taking expectations there, and using (53), we therefore obtain, for $\epsilon=\frac{1}{4 c_{\rho}^{2} T}$,

$$
\begin{align*}
\mathbb{E}\left[\Gamma_{t}\left|Y_{t}\right|^{2}\right] & +\frac{3}{4} \mathbb{E}\left[\int_{t}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right]+\mathbb{E}\left[\int_{t}^{T}\left(\kappa-\left(1+8 c_{f}^{2}+4 c_{f}^{2} c_{\rho}^{2}\right)\right) \Gamma_{s}\left|Y_{s}\right|^{2} d s\right] \\
\leq & \mathbb{E}\left[\Gamma_{T}\left|Y_{T}\right|^{2}\right]+\mathbb{E}\left[\int_{t}^{T} \Gamma_{s}\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} d s\right] \\
& +\frac{1}{4} \mathbb{E}\left[\sup _{0 \leq u \leq T} \Gamma_{u}\left|Y_{u}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right] \tag{54}
\end{align*}
$$

In particular, for $\kappa=1+8 c_{f}^{2}+4 c_{f}^{2} c_{\rho}^{2},(54)$ implies that

$$
\begin{align*}
& \sup _{t \in[0, T]} \mathbb{E}\left[\left|Y_{t}\right|^{2}\right]+\mathbb{E}\left[\int_{0}^{T}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right]  \tag{55}\\
& \quad \leq C \mathbb{E}\left(|\xi|^{2}+\int_{0}^{T}\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} d s+\sup _{0 \leq s \leq T}\left|Y_{s}\right|^{2}\right)
\end{align*}
$$

In order to show that the sup and the $\mathbb{E}$ can be switched on the left-hand side of (55), we rearrange the Itô formula (47) as

$$
\begin{align*}
\Gamma_{t}\left|Y_{t}\right|^{2}= & \Gamma_{T}|\xi|^{2}+2 \int_{t}^{T} \Gamma_{s}\left\langle Y_{s}, f\left(s, Y_{s}, Z_{s}, U_{s}, \rho_{s}(Y, Z, U)\right)\right\rangle d s-\kappa \int_{t}^{T} \Gamma_{s}\left|Y_{s}\right|^{2} d s \\
& -\int_{t}^{T} \Gamma_{s}\left|Z_{s}\right|^{2} d s-\int_{t}^{T} \int_{E} \Gamma_{s}\left|U_{s}(e)\right|^{2} j(d s, d e)  \tag{56}\\
& -2 \int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}-2 \int_{t}^{T} \Gamma_{s} \int_{E} Y_{s-} U_{s}(e) M(d s, d e) .
\end{align*}
$$

Plugging (50) and (52) into (56), we obtain

$$
\begin{align*}
\Gamma_{t}\left|Y_{t}\right|^{2} \leq & \Gamma_{T}|\xi|^{2}+\int_{t}^{T} \Gamma_{s}\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} d s+\frac{1}{4}\left[\int_{t}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right] \\
& +\left(1+8 c_{f}^{2}+4 c_{f}^{2} c_{\rho}^{2}-\kappa\right) \int_{t}^{T} \Gamma_{s}\left|Y_{s}\right|^{2} d s  \tag{57}\\
& +\epsilon c_{\rho}^{2} \int_{t}^{T} \Gamma_{s} \mathbb{E}_{s}\left(\sup _{s \leq u \leq T}\left|Y_{u}\right|^{2}+\int_{s}^{T}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u\right) d s \\
& -2 \int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}-2 \int_{t}^{T} \Gamma_{s} \int_{E} Y_{s-} U_{s}(e) M(d s, d e) \tag{58}
\end{align*}
$$

Setting $\kappa=1+8 c_{f}^{2}+4 c_{f}^{2} c_{\rho}^{2}$ yields

$$
\begin{align*}
\Gamma_{t}\left|Y_{t}\right|^{2} \leq & \Gamma_{T}|\xi|^{2}+\int_{t}^{T} \Gamma_{s}\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} d s+\frac{1}{4}\left[\int_{t}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right] \\
& +\epsilon c_{\rho}^{2} \int_{t}^{T} \Gamma_{s} \mathbb{E}_{s}\left(\sup _{s \leq u \leq T}\left|Y_{u}\right|^{2}+\int_{s}^{T}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u\right) g d s \\
& -2 \int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}-2 \int_{t}^{T} \Gamma_{s} \int_{E} Y_{s-} U_{s}(e) M(d s, d e) . \tag{59}
\end{align*}
$$

Taking the supremum over $t \in[0, T]$ and expectation on both sides, we have

$$
\begin{align*}
& \mathbb{E}\left(\sup _{t \in[0, T]} \Gamma_{t}\left|Y_{t}\right|^{2}\right) \leq \Gamma_{T} \mathbb{E}\left(|\xi|^{2}\right)+\mathbb{E}\left(\int_{0}^{T} \Gamma_{s}\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} d s\right) \\
& \quad+\frac{1}{4} \mathbb{E}\left(\int_{0}^{T} \Gamma_{s}\left(\left|Z_{s}\right|^{2}+\left|U_{s}\right|_{s}^{2}\right) d s\right) \\
& \quad+\epsilon c_{\rho}^{2} \mathbb{E}\left(\int_{0}^{T} \Gamma_{s} \mathbb{E}_{s}\left(\sup _{s \leq u \leq T}\left|Y_{u}\right|^{2}+\int_{s}^{T}\left(\left|Z_{u}\right|^{2}+\left|U_{u}\right|_{u}^{2}\right) d u\right) d s\right)  \tag{60}\\
& \quad+2 \mathbb{E}\left(\sup _{t \in[0, T]} \int_{t}^{T} \Gamma_{s} Y_{s}^{\top} Z_{s} d B_{s}\right)+2 \mathbb{E}\left(\sup _{t \in[0, T]} \int_{t}^{T} \Gamma_{s} \int_{E} Y_{s-} U_{s}(e) M(d s, d e)\right) .
\end{align*}
$$

Then, by using $\epsilon=\frac{1}{4 c_{\rho}^{2} T}$, (48), (49), (51), and the estimates on $Z$ and $U$ in (55), we deduce from (60)

$$
\mathbb{E}\left(\sup _{t \in[0, T]}\left|Y_{t}\right|^{2}\right) \leq C \mathbb{E}\left(|\xi|^{2}+\int_{0}^{T}\left|f\left(s, 0,0,0, \rho_{s}(0,0,0)\right)\right|^{2} d s\right)
$$

In conjunction with the estimates on $Z$ and $U$ in (55), this concludes the proof of (20).
A.2. Details of the proof of Theorem 3.1. Given $(X, V, W)$ and $\left(X^{\prime}, V^{\prime}, W^{\prime}\right)$ in $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times$ $\widehat{\mathcal{H}}_{2}^{l}$, let $(Y, Z, U)=\Phi(X, V, W)$ and $\left(Y^{\prime}, Z^{\prime}, U^{\prime}\right)=\Phi\left(X^{\prime}, V^{\prime}, W^{\prime}\right)$. For $L=X, V, W, Y, Z, U$, we denote $\delta L=L^{\prime}-L$. As in the proof of Lemma 3.1, letting $\Gamma_{t}=e^{\kappa t}$ and applying Ito's formula to $\Gamma|\delta Y|^{2}$ yields

$$
\begin{equation*}
\mathbb{E}\left[\Gamma_{t}\left|\delta Y_{t}\right|^{2}\right]+\mathbb{E}\left[\int_{t}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}+\kappa\left|\delta Y_{s}\right|^{2}\right) d s\right]=2 \int_{t}^{T} \mathbb{E}\left[\Gamma_{s}\left\langle\delta Y_{s}, \delta f_{s}\right\rangle\right] d s \tag{61}
\end{equation*}
$$

where

$$
\delta f_{s}=f\left(s, Y_{s}, Z_{s}, U_{s}, \rho_{s}(X, V, W)\right)-f\left(s, Y_{s}^{\prime}, Z_{s}^{\prime}, U_{s}^{\prime}, \rho_{s}\left(X^{\prime}, V^{\prime}, W^{\prime}\right)\right)
$$

Given Assumptions 3.1 on $\rho$ and 3.2 on $f$, applying Young's inequality with $c_{1}, c_{2} \neq 0$, we have

$$
\begin{align*}
& 2 \int_{t}^{T} \Gamma_{s}\left\langle\delta Y_{s}, \delta f_{s}\right\rangle d s \leq \int_{t}^{T}\left(2 c_{1}^{2} c_{f}+c_{2}^{2} c_{f}\right) \Gamma_{s}\left|\delta Y_{s}\right|^{2} d s \\
& \quad+\frac{c_{f}}{c_{1}^{2}} \int_{t}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s+\frac{c_{f}}{c_{2}^{2}} \int_{t}^{T} \Gamma_{s}\left|\rho_{s}(X, V, W)-\rho_{s}\left(X^{\prime}, V^{\prime}, W^{\prime}\right)\right|^{2} d s \\
& \leq  \tag{62}\\
& \quad \int_{t}^{T}\left(2 c_{1}^{2} c_{f}+c_{2}^{2} c_{f}\right) \Gamma_{s}\left|\delta Y_{s}\right|^{2} d s+\frac{c_{f}}{c_{1}^{2}} \int_{t}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s \\
& \quad+\frac{c_{f}}{c_{2}^{2}} c_{\rho}^{2} \int_{t}^{T} \Gamma_{s} \mathbb{E}_{s}\left[\sup _{s \leq u \leq T}\left|\delta X_{u}\right|^{2}+\int_{s}^{T}\left(\left|\delta V_{u}\right|^{2}+\left|\delta W_{u}\right|_{u}^{2}\right) d u\right] d s
\end{align*}
$$

Applying this inequality in (61) and proceeding as in (53) yields

$$
\begin{align*}
& \mathbb{E}\left[\Gamma_{t}\left|\delta Y_{t}\right|^{2}\right]+\mathbb{E}\left[\int_{t}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}+\kappa\left|\delta Y_{s}\right|^{2}\right) d s\right] \\
& \leq\left(2 c_{1}^{2} c_{f}+c_{2}^{2} c_{f}\right) \mathbb{E}\left[\int_{t}^{T} \Gamma_{s}\left|\delta Y_{s}\right|^{2} d s\right]+\frac{c_{f}}{c_{1}^{2}} \mathbb{E}\left[\int_{t}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s\right]  \tag{63}\\
& \quad+\frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}} \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right]
\end{align*}
$$

In particular, we obtain

$$
\begin{array}{rl}
\sup _{0 \leq t \leq T} & \mathbb{E}\left[\Gamma_{t}\left|\delta Y_{t}\right|^{2}\right]+\left(1-\frac{c_{f}}{c_{1}^{2}}\right) \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s\right] \\
+\left(\kappa-2 c_{1}^{2} c_{f}-c_{2}^{2} c_{f}\right) \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left|\delta Y_{s}\right|^{2} d s\right]  \tag{64}\\
\leq & \frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}} \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right]
\end{array}
$$

Choosing $c_{1}^{2}>c_{f}$ and $\kappa \geq 2 c_{1}^{2} c_{f}+c_{2}^{2} c_{f}$, we deduce from (64) that

$$
\begin{align*}
& \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s\right]  \tag{65}\\
& \quad \leq \frac{c_{1}^{2}}{\left(c_{1}^{2}-c_{f}\right)} \frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}} \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right]
\end{align*}
$$

For estimating the $\delta Y$ term, we write the Itô formula that applies to $\Gamma|\delta Y|^{2}$ in the manner of (56), i.e.,

$$
\begin{aligned}
\Gamma_{t}\left|\delta Y_{t}\right|^{2}= & 2 \int_{t}^{T} \Gamma_{s}\left\langle\delta Y_{s}, \delta f_{s}\right\rangle d s-\kappa \int_{t}^{T} \Gamma_{s}\left|\delta Y_{s}\right|^{2} d s \\
& -\int_{t}^{T} \Gamma_{s}\left|\delta Z_{s}\right|^{2} d s-\int_{t}^{T} \int_{E} \Gamma_{s}\left|\delta U_{s}(e)\right|^{2} j(d s, d e) \\
& -2 \int_{t}^{T} \Gamma_{s} \delta Y_{s}^{\top} \delta Z_{s} d B_{s}-2 \int_{t}^{T} \Gamma_{s} \int_{E} \delta Y_{s-} \delta U_{s}(e) M(d s, d e)
\end{aligned}
$$

Taking the supremum over $t \in[0, T]$ and expectation on both sides, using (62) and proceeding as in (53), we obtain

$$
\begin{align*}
& \mathbb{E}\left[\sup _{0 \leq t \leq T} \Gamma_{t}\left|\delta Y_{t}\right|^{2}\right]+\left(\kappa-2 c_{1}^{2} c_{f}-c_{2}^{2} c_{f}\right) \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left|\delta Y_{s}\right|^{2} d s\right] \leq \frac{c_{f}}{c_{1}^{2}} \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s\right]  \tag{66}\\
&+\frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}} \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right] \\
&+2 \mathbb{E}\left[\sup _{0 \leq t \leq T}\left|\int_{t}^{T} \Gamma_{s} \delta Y_{s}^{\top} \delta Z_{s} d B_{s}\right|\right]+2 \mathbb{E}\left[\sup _{0 \leq t \leq T}\left|\int_{t}^{T} \Gamma_{s} \int_{E} \delta Y_{s} \delta U_{s}(e) M(d s, d e)\right|\right]
\end{align*}
$$

By the Burkholder-Davis-Gundy inequality, there exists a positive constant $\bar{c}$ such that

$$
\begin{aligned}
\mathbb{E}\left[\sup _{t \in[0, T]}\left|\int_{t}^{T} \Gamma_{s} \delta Y_{s}^{\top} \delta Z_{s} d B_{s}\right|\right] & \leq \bar{c} \mathbb{E}\left[\left(\int_{0}^{T} e^{2 \kappa s}\left|Y_{s}\right|^{2}\left|Z_{s}\right|^{2} d s\right)^{1 / 2}\right] \\
& \leq \bar{c} \mathbb{E}\left[\left(\sup _{t \in[0, T]} \Gamma_{t}\left|Y_{t}\right|^{2}\right)^{1 / 2}\left(\int_{0}^{T} \Gamma_{s}\left|Z_{s}\right|^{2} d s\right)^{1 / 2}\right]
\end{aligned}
$$

Therefore, by Young's inequality with $c_{3}>0$ and (65) we have

$$
\begin{aligned}
& 2 \mathbb{E}\left[\sup _{t \in[0, T]}\left|\int_{t}^{T} \Gamma_{s} \delta Y_{s}^{\top} \delta Z_{s} d B_{s}\right|\right] \leq \frac{\bar{c}}{c_{3}^{2}} \mathbb{E}\left(\sup _{t \in[0, T]} \Gamma_{s}\left|\delta Y_{t}\right|^{2}\right)+c_{3}^{2} \bar{c} \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left|Z_{s}\right|^{2} d s\right] \\
& \leq \frac{\bar{c}}{c_{3}^{2}} \mathbb{E}\left(\sup _{t \in[0, T]} \Gamma_{s}\left|\delta Y_{t}\right|^{2}\right) \\
& \quad+c_{3}^{2} \bar{c} \frac{c_{1}^{2}}{\left(c_{1}^{2}-c_{f}\right)} \frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}} \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right] .
\end{aligned}
$$

Similarly, we can bound the last term in (66) as follows:

$$
\begin{align*}
& 2 \mathbb{E}\left[\sup _{0 \leq t \leq T}\left|\int_{t}^{T} \Gamma_{s} \int_{E} \delta Y_{s-} \delta U_{s}(e) M(d s, d e)\right|\right] \leq \frac{\bar{c}}{c_{3}^{2}} \mathbb{E}\left(\sup _{t \in[0, T]} \Gamma_{s}\left|\delta Y_{t}\right|^{2}\right) \\
& \quad+c_{3}^{2} \bar{c} \frac{c_{1}^{2}}{\left(c_{1}^{2}-c_{f}\right)} \frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}} \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right] . \tag{68}
\end{align*}
$$

Combining inequalities (65) through (68) gives

$$
\begin{aligned}
& \left(1-\frac{2 \bar{c}}{c_{3}^{2}}\right) \mathbb{E}\left[\sup _{0 \leq t \leq T} \Gamma_{t}\left|\delta Y_{t}\right|^{2}\right]+\left(1-\frac{c_{f}}{c_{1}^{2}}\right) \mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s\right] \\
& \leq\left(\frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}}+\left(1+2 c_{3}^{2} \bar{c}\right) \frac{c_{1}^{2}}{\left(c_{1}^{2}-c_{f}\right)} \frac{c_{f} c_{\rho}^{2} T}{c_{2}^{2}}\right) \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right] .
\end{aligned}
$$

Finally, we can choose $c_{1}^{2}=2 c_{f}, c_{3}^{2}=4 \bar{c}$, fix $c_{2}$ large enough so that $K=\frac{2\left(3+16 \bar{c}^{2}\right) c_{f} c_{\rho}^{2} T}{c_{2}^{2}}<1$, and we then have, for $\kappa \geq 2 c_{1}^{2} c_{f}+c_{2}^{2} c_{f}$ in $\Gamma=e^{\kappa \cdot}$,

$$
\begin{aligned}
& \mathbb{E}\left[\sup _{0 \leq t \leq T} \Gamma_{t}\left|\delta Y_{t}\right|^{2}\right]+\mathbb{E}\left[\int_{0}^{T} \Gamma_{s}\left(\left|\delta Z_{s}\right|^{2}+\left|\delta U_{s}\right|_{s}^{2}\right) d s\right] \\
& \quad \leq K \mathbb{E}\left[\sup _{0 \leq s \leq T} \Gamma_{s}\left|\delta X_{s}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|\delta V_{s}\right|^{2}+\left|\delta W_{s}\right|_{s}^{2}\right) d s\right] .
\end{aligned}
$$

Hence $\Phi$ is a contraction on the Banach space $\mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$ endowed with the norm defined by the square root of

$$
\mathbb{E}\left[\sup _{0 \leq t \leq T} \Gamma_{t}\left|X_{t}\right|^{2}+\int_{0}^{T} \Gamma_{s}\left(\left|V_{s}\right|^{2}+\left|W_{s}\right|_{s}^{2}\right) d s\right]
$$

for any $(X, V, W) \in \mathcal{S}_{2}^{l} \times \mathcal{H}_{2}^{l} \times \widehat{\mathcal{H}}_{2}^{l}$.
In order to prove the geometrical convergence of the Picard iteration, we can apply the previous computations with the following correspondence:

$$
\begin{array}{ll}
\delta Y=Y^{(n+1)}-Y^{(n)}, & \delta X=Y^{(n)}-Y^{(n-1)}, \\
\delta Z=Z^{(n+1)}-Z^{(n)}, & \delta V=Z^{(n)}-Z^{(n-1)}, \\
\delta U=U^{(n+1)}-U^{(n)}, & \delta W=U^{(n)}-U^{(n-1)},
\end{array}
$$

which shows that

$$
\begin{equation*}
\left\|Y^{(n+1)}-Y^{(n)}\right\|_{\mathcal{S}_{2}^{l}}^{2}+\left\|Z^{(n+1)}-Z^{(n)}\right\|_{\mathcal{H}_{2}^{l}}^{2}+\left\|U^{(n+1)}-U^{(n)}\right\|_{\widehat{\mathcal{H}}_{2}^{l}}^{2} \leq C K^{n} \tag{69}
\end{equation*}
$$

Appendix B. Invariance times. For completeness we recall the results of Crépey and Song (2017) that are used in the proofs of Theorem 5.1 and Lemma 5.1.

We consider the enlargement of the filtration setup corresponding to the condition (B) in section 5.1, which we reproduce here for the reader's convenience.

Condition (B). $\forall t \geq 0$ and $B \in \mathfrak{G}_{t}, \exists B^{\prime} \in \mathfrak{F}_{t}$ such that $B \cap\{t<\tau\}=B^{\prime} \cap\{t<\tau\}$.
The $\mathbb{F}$ Azéma supermartingale $S$ of $\tau$, with $\mathbb{F}$ predictable projection ${ }^{\nu} S$ (see Dellacherie and Meyer (1975)), admits the canonical Doob-Meyer decomposition $S=Q-D$, where $Q$ (with $Q_{0}=S_{0}$ ) and $D\left(\right.$ with $\left.D_{0}=0\right)$ are the $(\mathbb{F}, \mathbb{Q})$ local martingale component and the $(\mathbb{F}, \mathbb{Q})$ drift of the $(\mathbb{F}, \mathbb{Q})$ supermartingale $S$.

We work with $\mathbb{F}$ semimartingales and $(\mathbb{F}, \mathbb{P})$ and $(\mathbb{F}, \mathbb{Q})$ local martingales on a predictable set of interval type $\mathcal{I}$ as defined in He, Wang, and Yan (1992, section VIII.3). We recall that, for any semimartingale $X$ and for any predictable $X$ integrable process $L, X$ is a local martingale on $\mathcal{I}$ (respectively, $Y=L \cdot X$ on $\mathcal{I}$ ) means that

$$
\begin{equation*}
\left.X^{\iota_{n}} \text { is a local martingale (respectively, } Y^{\iota_{n}}=L \cdot\left(X^{\iota_{n}}\right)\right) \tag{70}
\end{equation*}
$$

holds for at least one or, equivalently, any nondecreasing sequence of stopping times $\left(\iota_{n}\right)_{n \geq 0}$ such that $\cup\left[0, \iota_{n}\right]=\mathcal{I}$.

The stochastic exponential of a semimartingale $X$ is denoted by $\mathcal{E}(X)$.
Lemmas 2.2 and A. 1 in Crépey and Song (2017). Under the condition (B) we have the following.

1. Let $M$ be a $(\mathbb{G}, \mathbb{Q})$ local martingale stopped before $\tau$. For any $\mathbb{F}$ optional reduction $M^{\prime}$ of $M, M^{\prime}$ is an $\mathbb{F}$ semimartingale on $\left\{\mathrm{S}_{-}>0\right\}$ and

$$
\begin{equation*}
\mathrm{S}_{-} \cdot M^{\prime}+\left[\mathrm{S}, M^{\prime}\right] \text { is an }(\mathbb{F}, \mathbb{Q}) \text { local martingale on }\left\{\mathrm{S}_{-}>0\right\} . \tag{71}
\end{equation*}
$$

Conversely, for any $\mathbb{F}$ semimartingale $K$ on $\left\{\mathrm{S}_{-}>0\right\}$ such that $\mathrm{S}_{-} . K+[\mathrm{S}, K]$ is an $(\mathbb{F}, \mathbb{Q})$ local martingale on $\left\{\mathrm{S}_{-}>0\right\}$, $K^{\tau-}$ is a $(\mathbb{G}, \mathbb{Q})$ local martingale on $\mathbb{R}_{+}$.
2. The Azéma supermartingale S of $\tau$ admits the multiplicative decomposition

$$
\begin{equation*}
S=S_{0} \mathcal{Q D} \text { on }\left\{{ }^{p} S>0\right\}, \tag{72}
\end{equation*}
$$

where $\mathcal{Q}=\mathcal{E}\left(\frac{1}{\overline{\mathcal{S}}} \cdot \mathbb{Q}\right)$ is an $(\mathbb{F}, \mathbb{Q})$ local martingale on $\left\{{ }^{p S}>0\right\}$ and $\mathcal{D}=\mathcal{E}\left(-\frac{1}{S_{-}} \cdot \mathrm{D}\right)$ is an $\mathbb{F}$ predictable nonincreasing process on $\left\{\begin{array}{l} \\ \\ S\end{array}>0\right\}$. Moreover, if $\tau$ has a $(\mathbb{G}, \mathbb{Q})$ intensity $\gamma$, then $\mathbf{D}$ is continuous and

$$
\begin{equation*}
\mathcal{E}\left(-\frac{1}{\mathrm{~S}_{-}} . \mathrm{D}\right)=e^{-\int_{0} \gamma_{t}^{\prime} d t} \tag{73}
\end{equation*}
$$

holds on $\left\{\mathrm{S}_{-}>0\right\}$.
Likewise, we reproduce here the condition (A) from section 5.1.
Condition (A). For any ( $\mathbb{F}, \mathbb{P}$ ) local martingale $P$ on $[0, T], P^{\tau-}$ is a $(\mathbb{G}, \mathbb{Q})$ local martingale on $[0, T]$.

Theorem 3.2 in Crépey and Song (2017). The condition (A) holds if and only if the $(\mathbb{F}, \mathbb{Q})$ density process of $\mathbb{P}$ coincides with

$$
\begin{equation*}
\mathcal{E}\left(\mathbb{1}_{\{\mathcal{P S}>0\}} \frac{1}{\mu \mathcal{S}} \cdot Q\right)_{\cdot \wedge T} \tag{74}
\end{equation*}
$$

on $\left\{{ }^{P} S>0\right\} \cap[0, T]$.

Theorem 3.7 in Crépey and Song (2017). If the condition (A) holds and $\tau$ has a $(\mathbb{G}, \mathbb{Q})$ intensity, then

$$
\begin{equation*}
\left\{S_{-}>0\right\}=\left\{{ }^{n} S>0\right\}=\{S>0\} . \tag{75}
\end{equation*}
$$

In addition,

$$
\begin{align*}
& \text { a process } P \text { is an }(\mathbb{F}, \mathbb{P}) \text { local martingale on }\left\{\mathrm{S}_{-}>0\right\} \cap[0, T] \\
& \quad \text { if and only if }  \tag{76}\\
& \mathrm{S}_{-} \cdot P+[\mathrm{S}, P] \text { is an }(\mathbb{F}, \mathbb{Q}) \text { local martingale on }\left\{\mathrm{S}_{-}>0\right\} \cap[0, T] .
\end{align*}
$$

Acknowledgment. We are thankful to Romuald Élie for useful discussions.

## REFERENCES

A. Agarwal, S. D. Marco, E. Gobet, J. López-Salas, F. Noubiagain, and A. Zhou (2019), Numerical approximations of McKean anticipative backward stochastic differential equations arising in initial margin requirements, ESAIM Proc. Surveys, 65, pp. 1-26.
C. Albanese and L. Andersen (2015), FVA: What's wrong, and how to fix it, Risk Mag., September, pp. 5456.
C. Albanese, S. Caenazzo, and S. Crépey (2016), Capital and funding, Risk Mag., May, pp. 71-76.
C. Albanese, S. Caenazzo, and S. Crépey (2017), Credit, funding, margin, and capital valuation adjustments for bilateral portfolios, Probab. Uncertain. Quant. Risk, 2, 7.
C. Albanese and S. Crépey (2019), Capital Valuation Adjustment and Funding Valuation Adjustment, https: //math.maths.univ-evry.fr/crepey.
C. Albanese, S. Crépey, R. Hoskinson, and B. Saadeddine (2019), XVA Analysis from the Balance Sheet, https://math.maths.univ-evry.fr/crepey.
L. Andersen, D. Duffie, and Y. Song (2019), Funding value adjustments, J. Finance, 74, pp. 145-192.
P. Artzner, F. Delbean, J. Eber, and D. Heath (1999), Coherent measures of risk, Math. Finance, 9, pp. 203-228.
D. Barrera, S. Crépey, B. Diallo, G. Fort, E. Gobet, and U. Stazhynski (2019), Stochastic approximation schemes for economic capital and risk margin computations, ESAIM Proc. Surveys, 65, pp. 182-218.
Basel Committee on Banking Supervision (2011), Basel III: A Global Regulatory Framework for more Resilient Banks and Banking Systems, http://www.bis.org/publ/bcbs189.pdf.
M. Bichuch, A. Capponi, and S. Sturm (2018), Arbitrage-free XVA, Math. Finance, 28, pp. 582-620.
T. Bielecki and M. Rutkowski (2002), Credit Risk: Modeling, Valuation and Hedging, Springer Finance, Berlin.
T. R. Bielecki, M. Jeanblanc, and M. Rutkowski (2009), Credit Risk Modeling, CSFI Lecture Notes Ser. 2, Osaka University Press, Osaka University.
T. R. Bielecki and M. Rutkowski (2015), Valuation and hedging of contracts with funding costs and collateralization, SIAM J. Financial Math., 6, pp. 594-655.
D. Brigo and A. Capponi (2010), Bilateral counterparty risk with application to CDSs, Risk Mag., March, pp. 85-90.
D. Brigo and A. Pallavicini (2014), Nonlinear consistent valuation of CCP cleared or CSA bilateral trades with initial margins under credit, funding and wrong-way risks, J. Financ. Eng., 1, pp. 1-60.
C. Burgard and M. Kjaer (2011), In the balance, Risk Mag., October, pp. 72-75.
C. Burgard and M. Kjaer (2013), Funding costs, funding strategies, Risk Mag., December, pp. 82-87.
C. Burgard and M. Kjaer (2017), Derivatives funding, netting and accounting, Risk Mag., March, pp. 100104.
S. Crépey and S. Song (2017), Invariance times, Ann. Probab., 45, pp. 4632-4674.
C. Dellacherie and P.-A. Meyer (1975), Probabilité et Potentiel, Hermann, Paris.
D. Duffie and M. Huang (1996), Swap rates and credit quality, J. Finance, 51, pp. 921-950.
R. Elliot and Z. Yang (2013), Anticipated backward stochastic differential equations with continuous coefficients, Commun. Stoch. Anal., 7, pp. 303-319.
Y. Elouerkhaoui (2016), From FVA to KVA: including cost of capital in derivatives pricing, Risk Mag., March, pp. 70-75.
A. Green, C. Kenyon, and C. Dennis (2014), KVA: capital valuation adjustment by replication, Risk Mag., December, pp. 82-87.
S.-W. He, J.-G. Wang, and J.-A. Yan (1992), Semimartingale Theory and Stochastic Calculus, CRC, Coca Raton, FL.
J. Hull and A. White (2012A), The FVA debate, followed by The FVA debate continued, Risk Mag., July, pp. 83-85.
J. Hull and A. White (2012B), The FVA debate, followed by The FVA debate continued, Risk Mag., October, p. 52.
J. Jacod (1979), Calcul Stochastique et Problèmes de Martingales, Lecture Notes in Math. 714, Berlin, Springer.
T. Kruse and A. Popier (2016), BSDEs with monotone generator driven by Brownian and Poisson noises in a general filtration, Stochastics, 88, pp. 491-539.
A. Nikeghbali and M. Yor (2005), A definition and some characteristic properties of pseudo-stopping times, Ann. Probab., 33, pp. 1804-1824.
S. Peng and Z. Yang (2009), Anticipated backward stochastic differential equations, Ann. Probab., 37, pp. 877902.
M. Pykhtin, ed. (2005), Counterparty Credit Risk Modelling: Risk Management, Pricing and Regulation, Risk Books, London.
P. Schönbucher (2004), A measure of survival, Risk Mag., 17, pp. 79-85.


[^0]:    *Received by the editors February 4, 2019; accepted for publication (in revised form) November 19, 2019; published electronically February 18, 2020.
    https://doi.org/10.1137/19M1242781
    Funding: The research of Stéphane Crépey has been conducted with the support of the Research Initiative "Modélisation des marchés actions, obligations et dérivés," financed by HSBC France under the aegis of the Europlace Institute of Finance. The views and opinions expressed in this paper are those of the authors alone and do not necessarily reflect the views or policies of HSBC Investment Bank, its subsidiaries, or affiliates. The research of Wissal Sabbagh has been conducted with the support of the Fondation du Risque, program "XVA analysis and applications," under the aegis of Louis Bachelier laboratory and of the ANR project 11-LABX-0019.
    ${ }^{\dagger}$ LaMME, Univ Evry, CNRS, Université Paris-Saclay, 91037, Evry, France (stephane.crepey@univ-evry.fr, wissal.sabbagh@univ-evry.fr, shiqi.song@univ-evry.fr).

[^1]:    ${ }^{1}$ See https://www.risk.net/derivatives/2322843/jp-morgan-takes-15-billion-fva-loss.

[^2]:    ${ }^{2}$ For which $\sigma$ integrability of $X$ valued at any stopping time, e.g., $X$ bounded or càdlàg, is enough.

[^3]:    ${ }^{3}$ At least, the part of their results derived under square integrable assumptions, including their Theorem 1, which we use in the paper.

[^4]:    ${ }^{4}$ See Lemmas 5.2-5.3 in Albanese and Crépey (2019) for detailed derivations.

[^5]:    ${ }^{5}$ All the Lebesgue integrals that appear in these expressions are well defined over $[0, T]$, having assumed semimartingale XVAs.

[^6]:    ${ }^{6}$ For which having $\Lambda$ càdlàg is enough.

[^7]:    ${ }^{7}$ Here we mean square integrable solution in the sense immediately analogous to the one of Definition 4.1, as justified at the end of the proof of Theorem 4.1.

[^8]:    ${ }^{8}$ cf. Crépey and Song (2017, eq. (2.1)).
    ${ }^{9}$ Here we assume that the involved left-limit exists and we use the convention that $Y_{0-}=Y_{0}$.

[^9]:    ${ }^{10}$ The notation in the present paper differs from the one in Albanese and Crépey (2019). In particular, $\widetilde{\mathbb{E}}$ and $\mathbb{E}$ here correspond to $\mathbb{E}$ and $\mathbb{E}^{\prime}$ there; $\widetilde{\mathcal{S}}_{2}$ and $\mathcal{S}_{2}$ here correspond to $\mathcal{S}_{2}$ and $\mathcal{S}_{2}^{\prime}$ there; " $\sim$ " for " $(\mathbb{G}, \mathbb{Q})$ equations" is only used here; the risk-free asset is used as numéraire there.

