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# STRATIFIED REGRESSION MONTE-CARLO SCHEME FOR SEMILINEAR PDES AND BSDES WITH LARGE SCALE PARALLELIZATION ON GPUS 

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#### Abstract

In this paper, we design a novel algorithm based on Least-Squares Monte Carlo (LSMC) in order to approximate the solution of discrete time Backward Stochastic Differential Equations (BSDEs). Our algorithm allows massive parallelization of the computations on multicore devices such as graphics processing units (GPUs). Our approach consists of a novel method of stratification which appears to be crucial for large scale parallelization.


Key words. Backward stochastic differential equations, dynamic programming equation, empirical regressions, parallel computing, GPUs, CUDA.

AMS subject classifications. 49L20, 62Jxx, 65C30, 93E24, 68W10.

## 1. Introduction.

The problem. The aim of the algorithm in this paper is to approximate the ( $Y, Z$ ) components of the solution to the decoupled forward-backward stochastic differential equation (BSDE)

$$
\begin{align*}
& Y_{t}=g\left(X_{T}\right)+\int_{t}^{T} f\left(s, X_{s}, Y_{s}, Z_{s}\right) \mathrm{d} s-\int_{t}^{T} Z_{s} \mathrm{~d} W_{s},  \tag{1.1}\\
& X_{t}=x+\int_{0}^{t} b\left(s, X_{s}\right) \mathrm{d} s+\int_{0}^{t} \sigma\left(s, X_{s}\right) \mathrm{d} W_{s}, \tag{1.2}
\end{align*}
$$

where $W$ is a $q \geq 1$ dimensional Brownian motion. The algorithm will also approximate the solution $u$ to the related semilinear, parabolic partial differential equation (PDE) of the form

$$
\begin{equation*}
\partial_{t} u(t, x)+\mathcal{A} u(t, x)+f\left(t, x, u(t, x), \nabla_{x} u \sigma(t, x)\right)=0 \text { for } t<T \text { and } u(T, .)=g(.), \tag{1.3}
\end{equation*}
$$

where $\mathcal{A}$ is the infinitesimal generator of $X$, through the Feynman-Kac relation $\left(Y_{t}, Z_{t}\right)=\left(u\left(t, X_{t}\right),\left(\nabla_{x} u \sigma\right)\left(t, X_{t}\right)\right)$. In recent times, there has been an increasing interest to have algorithms which work efficiently when the dimension $d$ of the space occupied by the process $X$ is large. This interest has been principally driven by the mathematical finance community, where nonlinear valuation rules are becoming increasingly important.

In general, currently available algorithms $[6,3,4,14,5,10,11,9]$ rarely handle the case of dimension greater than 8 . The main constraint is not only due to the computational time, but mainly due to memory consumption requirements by the algorithms. For example, the recent work [11] uses a Regression Monte Carlo approach (a.k.a. Least Squares MC), in which the solutions ( $u, \nabla_{x} u \sigma$ ) of the semi-linear PDE

[^0]are approximated on a $\mathcal{K}$-dimensional basis of functions at each point of a time grid of cardinality $N$. The coefficients of the basis functions are computed at every time point $t_{i}$ with the aid of $\mathcal{M}$ simulations of a discrete time Markov chain (which approximates $X$ ) in the interval $\left[t_{i}, T\right]$. The main memory constraints of this scheme are (a) to store the $\mathcal{K} \times N$ coefficients of the basis functions, and (b) to store the $\mathcal{M} \times N$ simulations used to compute the coefficients. To illustrate the problem of high dimension, the dimension of the basis is typically $\mathcal{K}=$ const $\times N^{\alpha d}$, for some $\alpha>0$ (which decreases with the regularity of the solution), so $\mathcal{K}$ increases geometrically with $d$. Moreover, the error analysis of these algorithms demonstrates that the local statistical error is proportional to $N \mathcal{K} / \mathcal{M}$, so that one must choose $\mathcal{M}=$ const $\times \mathcal{K} N^{2}$ to ensure a convergence $O\left(N^{-1}\right)$ of the scheme. This implies that the simulations pose by far the most significant constraint on the memory.

Objectives. The purpose of this paper is to drastically rework the algorithm of [11] to first minimize the exposure to the memory due to the storage of simulations. This will allow computation in larger dimension $d$. Secondly, in this way the algorithm can be implemented in parallel on GPU processors to optimize the computational time.

New Regression Monte Carlo paradigm. We develop a novel algorithm called the Stratified Regression MDP (SRMDP) algorithm; the name is aimed to distinguish from the related LSMDP algorithm [11]. The key technique is to use stratified simulation of the paths of $X$. In order to estimate the solution at $t_{i}$, we first define a set of hypercubes $\left(\mathcal{H}_{k} \subset \mathbb{R}^{d}: 1 \leq k \leq K\right)$. Then, for each hypercube $\mathcal{H}_{k}$, we simulate $M$ paths of the process $X$ in the interval $\left[t_{i}, T\right]$ starting from i.i.d. random variables valued in $\mathcal{H}_{k}$; these random variables are distributed according to the conditional logistic distribution, see $\left(\mathbf{A}_{\nu}\right)$ later. By using only the paths starting in $\mathcal{H}_{k}$, we approximate the solution to the BSDE restricted to $X_{t_{i}} \in \mathcal{H}_{k}$ on linear functions spaces $\mathcal{L}_{Y, k}$ and $\mathcal{L}_{Z, k}$ (both of small dimension), see ( $\mathbf{A}_{\text {Strat. }}$ ) later. This allows us to minimize the amount of memory consumed by the simulations, since we only need to generate samples on one hypercube at a time. In Theorem 3.5, we demonstrate that the error of our scheme is proportional to $N \max \left(\operatorname{dim}\left(\mathcal{L}_{Y, k}\right), \operatorname{dim}\left(\mathcal{L}_{Z, k}\right)\right) / M$ and, since $\max \left(\operatorname{dim}\left(\mathcal{L}_{Y, k}\right), \operatorname{dim}\left(\mathcal{L}_{Z, k}\right)\right)=$ const, we require only $M=$ const $\times N^{2}$ to ensure the convergence $O\left(N^{-1}\right)$. Therefore, the memory consumption of the algorithm will be dominated by the storage of the coefficients, which equals const $\times N^{-\alpha d}$ (the theoretical minimum). Moreover, the computations are performed in parallel across the hypercubes, which allows for massive parallelization. The speed-up compared to sequential programming increases as the dimension $d$ increases, because of the geometric growth of the number of hypercubes with respect to $d$. In the subsequent tests ( $\S 5$ ), for instance we can solve problems in dimension $d=11$ within a couple of seconds using 100 simulations per hypercube.

This regression Monte Carlo approach is very different from the algorithm proposed in [11]. Although local approximations were already proposed in that work, the paths of the process $X$ were simulated from a fixed point at time 0 rather than directly in the hypercubes. This implies that one must store all the simulated paths at any given time, rather than only those for the specific hypercubes. This is because the trajectories are random, and one is not certain which paths will end up in which hypercubes a priori. Therefore, our scheme essentially removes the main constraint on the memory consumption of LSMC algorithms for BSDEs.

The choice of the logistic distribution for the stratification procedure is crucial. Firstly, it is easy to simulate from the conditional distribution. Secondly, it possesses the important USES property (see later $\left(\mathbf{A}_{\nu}\right)$ ), which enables us to recover equivalent
$\mathbf{L}_{2}$-norms (up to constant) for the marginal of the forward process initialized with the logistic distribution (Proposition 2.1).

Literature review. Parallelization of Monte-Carlo methods for solving non-linear probabilistic equations has been shortly investigated. Due to the non-linearity, this is a challenging issue. For optimal stopping problems, we can refer to the works [1] and [2] with numerical results in dimension 4. To the best of our knowledge, the only work related to BSDEs in parallel version is [14]. It is based on a Picard iteration for finding the solution, coupled with iterative control variates. The iterative solution is computed through an approximation on sparse polynomial basis. Although the authors report efficient numerical experiments up to dimension 8 , this study is not supported by a theoretical error analysis. Due to the stratification, our proposed approach is quite different from [14] and additionally, we provide an error analysis (Theorem 3.5).

## Notation.

(i) $|x|$ stands for the Euclidean norm of the vector $x$.
(ii) $\log (x)$ stands for the natural logarithm of $x \in \mathbb{R}_{+}$.
(iii) For a multidimensional process $U=\left(U_{i}\right)_{0 \leq i \leq N}$, its $l$-th component is denoted by $U_{l}=\left(U_{l, i}\right)_{0 \leq i \leq N}$.
(iv) For any finite $L>0$ and $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$, define the truncation function

$$
\begin{equation*}
\mathcal{T}_{L}(x):=\left(-L \vee x_{1} \wedge L, \ldots,-L \vee x_{n} \wedge L\right) \tag{1.4}
\end{equation*}
$$

(v) For a probability measure $\nu$ on a domain $D$, and function $h: D \rightarrow \mathbb{R}^{l}$ in $\mathbf{L}_{2}(D, \nu)$, denote the $\mathbf{L}_{2}$ norm of $h$ by $|h|_{\nu}:=\sqrt{\int_{D}|h|^{2}(x) \nu(\mathrm{d} x)}$.
(vi) For a probability measure $\nu$, disjoint sets $\left\{\mathcal{H}_{1}, \ldots, \mathcal{H}_{K}\right\}$ in the support of $\nu$, and finite dimensional function spaces $\mathcal{L}\left\{\mathcal{L}_{1}, \ldots, \mathcal{L}_{K}\right\}$ such that the domain of $\mathcal{L}_{k}$ is in the respected set $\mathcal{H}_{k}$

$$
\nu(\operatorname{dim}(\mathcal{L}))=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) \operatorname{dim}\left(\mathcal{L}_{k}\right)
$$

(vii) For function $g: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$, the order notation $g(x)=O(x)$ means that there exists some universal unspecified constant, const $>0$, such that $g(x) \leq$ const $\times x$ for all $x \in \mathbb{R}_{+}$.
2. Mathematical framework and basic properties. We work on a filtered probability space $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{0 \leq t \leq T}, \mathbb{P}\right)$ containing a $q$-dimensional $(q \geq 1)$ Brownian motion $W$. The filtration $\left(\mathcal{F}_{t}\right)_{0 \leq t \leq T}$ satisfies the usual hypotheses. The existence of a unique strong solution $X$ to the forward equation (1.2) follows from usual Lipschitz conditions on $b$ and $\sigma$, see ( $\mathbf{A}_{\mathbf{X}}$ ). The BSDE (1.1) is approximated using a multistepforward dynamical programming equation (MDP) studied in [10]. Let $\pi:=\left\{t_{i}:=\right.$ $\left.i \Delta_{t}: 0 \leq i \leq N\right\}$ be the uniform time-grid with time step $\Delta_{t}=T / N$. The solution $\left(Y_{i}, Z_{i}\right)_{0 \leq i \leq N-1}$ of the MDP can be written in the form:

$$
\left.\begin{array}{rl}
Y_{i} & =\mathbb{E}_{i}\left(g\left(X_{N}\right)+\sum_{j=i}^{N-1} f_{j}\left(X_{j}, Y_{j+1}, Z_{j}\right) \Delta_{t}\right)  \tag{2.1}\\
\Delta_{t} Z_{i} & =\mathbb{E}_{i}\left(\left(g\left(X_{N}\right)+\sum_{j=i+1}^{N-1} f_{j}\left(X_{j}, Y_{j+1}, Z_{j}\right) \Delta_{t}\right) \Delta W_{i}\right)
\end{array}\right\} \quad \text { for } i \in\{0, \ldots, N-1\}
$$

where $\left(X_{j}\right)_{i \leq j \leq N}$ is a Markov chain approximating the forward component (1.2) (typically the Euler scheme, see Algorithm 2 below), $\Delta W_{i}:=W_{t_{i+1}}-W_{t_{i}}$ is the
$(i+1)$-th Brownian motion increment, and $\mathbb{E}_{i}(\cdot):=\mathbb{E}\left(\cdot \mid \mathcal{F}_{t_{i}}\right)$ is the conditional expectation. Our working assumptions on the functions $g$ and $f$ are as follows:
$\left(\mathbf{A}_{\mathbf{g}}\right) g$ is a bounded measurable function from $\mathbb{R}^{d}$ to $\mathbb{R}$, the upper bound of which is denoted by $C_{g}$.
$\left(\mathbf{A}_{\mathbf{f}}\right)$ for every $i<N, f_{i}(x, y, z)$ is a measurable function $\mathbb{R}^{d} \times \mathbb{R} \times \mathbb{R}^{q}$ to $\mathbb{R}$, and there exist two finite constants $L_{f}$ and $C_{f}$ such that, for every $i<N$,

$$
\begin{aligned}
&\left|f_{i}(x, y, z)-f_{i}\left(x, y^{\prime}, z^{\prime}\right)\right| \leq L_{f}\left(\left|y-y^{\prime}\right|+\left|z-z^{\prime}\right|\right) \\
& \forall\left(x, y, y^{\prime}, z, z^{\prime}\right) \in \mathbb{R}^{d} \times(\mathbb{R})^{2} \times\left(\mathbb{R}^{q}\right)^{2} \\
&\left|f_{i}(x, 0,0)\right| \leq C_{f}, \quad \forall x \in \mathbb{R}^{d}
\end{aligned}
$$

The definition of the Markov chain $\left(X_{j}\right)_{j}$ is made under the following assumptions. $\left(\mathbf{A}_{\mathbf{X}}\right)$ The coefficients functions $b$ and $\sigma$ satisfy
(i) $b:[0, T] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ and $\sigma:[0, T] \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d} \otimes \mathbb{R}^{q}$ are bounded measurable, uniformly Lipschitz in the space dimensions;
(ii) there exists $\zeta \geq 1$ such that, for all $\xi \in \mathbb{R}^{d}$, the following inequalities hold: $\zeta^{-1}|\xi|^{2} \leq \xi^{\top} \sigma(t, x) \sigma(t, x)^{\top} \xi \leq \zeta|\xi|^{2}$.
Let $X_{i}$ be a random variable with some distribution $\eta$ (more details on this to follow). Then $X_{j}$ for $j>i$ is generated according to one of the two algorithms below:

Algorithm 1 (SDE dynamics). $\quad X_{j+1}=\bar{X}_{t_{j+1}}=X_{j}+\int_{t_{j}}^{t_{j+1}} b\left(s, \bar{X}_{s}\right) \mathrm{d} s+$ $\int_{t_{j}}^{t_{j+1}} \sigma\left(s, \bar{X}_{s}\right) \mathrm{d} W_{s} ;$

Algorithm 2 (Euler dynamics). $X_{j+1}=X_{j}+b\left(t_{i}, X_{i}\right) \Delta_{t}+\sigma\left(t_{i}, X_{i}\right) \Delta W_{i}$.

The above ellipticity condition (ii) will be used in the proof of Proposition 2.1.
As in the continuous time framework (1.1), the solution of the MDP (2.1) admits a Markov representation: under $\left(\mathbf{A}_{\mathbf{X}}\right)$, for every $i$, there exist measurable deterministic functions $y_{i}: \mathbb{R}^{d} \rightarrow \mathbb{R}$ and $z_{i}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{q}$, such that $Y_{i}=y_{i}\left(X_{i}\right)$ and $Z_{i}=z_{i}\left(X_{i}\right)$, almost surely. In fact, the value functions $y_{i}(\cdot)$ and $z_{i}(\cdot)$ are independent of how we initialize the forward component.

For the subsequent stratification algorithm, $X_{i}$ will be sampled randomly (and independently of the Brownian motion $W$ ) according to different squared-integrable distributions $\eta$. When $X_{i} \sim \eta$, we will write $\left(X_{j}^{(i, \eta)}\right)_{i \leq j \leq N}$ the Markov chain given in ( $\mathbf{A}_{\mathbf{X}}$ ), using either the SDE dynamics (better when possible) or the Euler one. One can recover the value functions from the conditional expectations: almost surely,

$$
\begin{equation*}
y_{i}\left(X_{i}^{(i, \eta)}\right)=\mathbb{E}\left(g\left(X_{N}^{(i, \eta)}\right)+\sum_{j=i}^{N-1} f_{j}\left(X_{j}^{(i, \eta)}, y_{j+1}\left(X_{j+1}^{(i, \eta)}\right), z_{j}\left(X_{j}^{(i, \eta)}\right)\right) \Delta_{t} \mid X_{i}^{(i, \eta)}\right) \tag{2.2}
\end{equation*}
$$

$$
\Delta_{t} z_{i}\left(X_{i}^{(i, \eta)}\right)=\mathbb{E}\left(\left(g\left(X_{N}^{(i, \eta)}\right)+\sum_{j=i+1}^{N-1} f_{j}\left(X_{j}^{(i, \eta)}, y_{j+1}\left(X_{j+1}^{(i, \eta)}\right), z_{j}\left(X_{j}^{(i, \eta)}\right)\right) \Delta_{t}\right) \Delta W_{i} \mid X_{i}^{(i, \eta)}\right)
$$

the proof of this is the same as [11, Lemma 4.1].
Approximating the solution to (2.1) is actually achieved by approximating the functions $y_{i}(\cdot)$ and $z_{i}(\cdot)$. In this way, we are directly approximating the solution to the semilinear PDE (1.3). Our approach consists in approximating the restrictions of the functions $y_{i}$ and $z_{i}$ to subsets of a cubic partition of $\mathbb{R}^{d}$ using finite dimensional linear function spaces. The basic assumptions for this local approximation approach are given below.
( $\mathbf{A}_{\text {Strat. }}$ ) There are $K \in \mathbb{N}^{*}$ disjoint hypercubes $\left(\mathcal{H}_{k}: 1 \leq k \leq K\right)$, that is

$$
\mathcal{H}_{k} \cap \mathcal{H}_{l}=\emptyset, \quad \bigcup_{k=1}^{K} \mathcal{H}_{k}=\mathbb{R}^{d} \quad \text { and } \quad \mathcal{H}_{k}=\prod_{l=1}^{d}\left[x_{k, l}^{-}, x_{k, l}^{+}\right)
$$

for some $-\infty \leq x_{k, l}^{-}<x_{k, l}^{+} \leq+\infty$. Additionally, there are linear function spaces $\mathcal{L}_{Y, k}$ and $\mathcal{L}_{Z, k}$, valued in $\mathbb{R}$ and $\mathbb{R}^{q}$ respectively, which are subspaces of $\mathbf{L}_{2}\left(H_{k}, \nu_{k}\right)$ w.r.t. a probability measure $\nu_{k}$ on $\mathcal{H}_{k}$ defined in $\left(\mathbf{A}_{\nu}\right)$ below. Common examples of hypercubes are:
(i) Hypercubes of equal size: $x_{k, l}^{+}-x_{k, l}^{-}=$const $>0$ for all $k$ and $l$, except for exterior strata that must be infinite.
(ii) Hypercubes of equal probability: $\nu\left(\mathcal{H}_{k}\right)=1 / K$ for some probability $\nu$ to be defined later in $\left(\mathbf{A}_{\nu}\right)$.

Common examples of local approximations spaces $\mathcal{L}_{Y, k}$ and $\mathcal{L}_{Z, k}$ are:
(i) Piece-wise constant approximation $(\mathbf{L P 0}): \mathcal{L}_{Y, k}:=\operatorname{span}\left\{\mathbf{1}_{\mathcal{H}_{k}}\right\}$, and $\mathcal{L}_{Z, k}:=$ $\left(\mathcal{L}_{Y, k}\right)^{q} ; \operatorname{dim}\left(\mathcal{L}_{Y}\right)=1$ and $\operatorname{dim}\left(\mathcal{L}_{Z, k}\right)=q$.
(ii) Affine approximations $(\mathbf{L P} 1): \mathcal{L}_{Y, k}:=\operatorname{span}\left\{\mathbf{1}_{\mathcal{H}_{k}}, x_{1} \mathbf{1}_{\mathcal{H}_{k}}, \ldots, x_{d} \mathbf{1}_{\mathcal{H}_{k}}\right\}$, and $\mathcal{L}_{Z, k}:=\left(\mathcal{L}_{Y, k}\right)^{q} ; \operatorname{dim}\left(\mathcal{L}_{Y}\right)=d+1$ and $\operatorname{dim}\left(\mathcal{L}_{Z, k}\right)=q(d+1)$.

The key idea in this paper is to select a distribution $\nu$, the restriction of which to the hypercubes $\mathcal{H}_{k}, \nu_{k}$, can be explicitly computed. Then, we can easily simulate i.i.d. copies of $X_{i}^{\left(i, \nu_{k}\right)}$ directly in $\mathcal{H}_{k}$ and use the resulting paths of the Markov chain to estimate $\left.y_{k}(\cdot)\right|_{\mathcal{H}_{k}}$. This sampling method is traditionally known as stratification, and for this reason we will call the hypercubes in $\left(\mathbf{A}_{\text {Strat. }}\right)$ the strata. For the stratification, the components $X_{i}^{\left(i, \nu_{k}\right)}$ are sampled as i.i.d. conditional logistic random variables, which is precisely stated in the following assumption.
$\left(\mathbf{A}_{\nu}\right)$ Let $\mu>0$. The distribution of $X_{i}^{\left(i, \nu_{k}\right)}$ is given by $\mathbb{P} \circ\left(X_{i}^{\left(i, \nu_{k}\right)}\right)^{-1}(\mathrm{~d} x)=\nu_{k}(\mathrm{~d} x)$, where

$$
\nu_{k}(\mathrm{~d} x)=\frac{\mathbf{1}_{\mathcal{H}_{k}}(x) \nu(\mathrm{d} x)}{\nu\left(\mathcal{H}_{k}\right)}
$$

and

$$
\nu(\mathrm{d} x)=p_{\text {logis. }}^{(\mu)}(x) \mathrm{d} x, \quad p_{\text {logis. }}^{(\mu)}(x):=\prod_{l=1}^{d} \frac{\mu e^{-\mu x_{l}}}{\left(1+e^{-\mu x_{l}}\right)^{2}}, \quad x=\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d}
$$

REMARK 2.1. An important relation of $\nu$ and $\nu_{k}$ is that one has the $\mathbf{L}_{2}$-norm identity $|\cdot|_{\nu}^{2}=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right)|\cdot|_{\nu_{k}}^{2}$.

In order to generate the random variable $X_{i}^{\left(i, \nu_{k}\right)}$, we make use of the inverse conditional distribution function of $\nu_{k}$ and the simulation of uniform random variables, as shown in the following algorithm:

Algorithm 3. Draw d independent random variables $\left(U_{1}, \ldots, U_{d}\right)$ which are uniformly distributed on $[0,1]$, and compute

$$
X_{i}^{\left(i, \nu_{k}\right)}:=\left(F_{\nu,\left[x_{k, 1}^{-}, x_{k, 1}^{+}\right)}^{-1}\left(U_{1}\right), \ldots, F_{\nu,\left[x_{k, d}^{-}, x_{k, d}^{+}\right)}^{-1}\left(U_{d}\right)\right) \stackrel{d}{\sim} \nu_{k}
$$

where we use the functions $F_{\nu}(x):=\int_{-\infty}^{x} \nu\left(\mathrm{~d} x^{\prime}\right)=1 /(1+\exp (-\mu x))$ and

$$
F_{\nu,\left[x^{-}, x^{+}\right)}^{-1}(U)=-\frac{1}{\mu} \log \left(\frac{1}{F_{\nu}\left(x^{-}\right)+U\left(F_{\nu}\left(x^{+}\right)-F_{\nu}\left(x^{-}\right)\right)}-1\right)
$$

A further reason for the choice of the logistic distribution is that it induces the following stability property on the $\mathbf{L}_{2}$ norms of the Markov chain $\left(X_{j}^{(i, \nu)}\right)_{i \leq j \leq N}$; this property will be crucial for the error analysis of the stratified regression scheme in §3.2. The proof is postponed to Appendix A.1.

Proposition 2.1. Suppose that $\nu$ is the logistic distribution defined in $\left(\mathbf{A}_{\nu}\right)$. There is a constant $c_{\left(\mathbf{A}_{\nu}\right)} \in[1,+\infty)$ such that, for any function $h: \mathbb{R}^{d} \mapsto \mathbb{R}$ or $\mathbb{R}^{q}$ in $\mathbf{L}_{2}(\nu)$, for any $0 \leq i \leq N$, and for any $i \leq j \leq N-1$, we have

$$
\frac{1}{c_{\left(\mathbf{A}_{\nu}\right)}} \mathbb{E}\left[\left|h\left(X_{j}^{(i, \nu)}\right)\right|^{2}\right] \leq|h|_{\nu}^{2} \leq c_{\left(\mathbf{A}_{\nu}\right)} \mathbb{E}\left[\left|h\left(X_{j}^{(i, \nu)}\right)\right|^{2}\right] .
$$

To conclude this section, we recall standard uniform absolute bounds for the functions $y_{i}(\cdot)$ and $z_{i}(\cdot)$.

Proposition 2.2 (a.s. upper bounds, [11, Proposition 3.3]). For $N$ large enough such that $\frac{T}{N} L_{f}^{2} \leq \frac{1}{12 q}$, we have for any $x \in \mathbb{R}^{d}$ and any $0 \leq i \leq N-1$,

$$
\begin{equation*}
\left|y_{i}(x)\right| \leq C_{y}:=e^{\frac{T}{4}+6 q\left(1 \vee L_{f}^{2}\right)(T \vee 1)}\left(C_{g}+\frac{T}{2 \sqrt{q}} C_{f}\right), \quad\left|z_{l, i}(x)\right| \leq C_{z}:=\frac{C_{y}}{\sqrt{\Delta_{t}}} \tag{2.3}
\end{equation*}
$$

## 3. Stratified algorithm and convergence results.

3.1. Algorithm. In this section, we define the SRMDP algorithm mathematically, and then expose in $\S 4$ how to efficiently perform it using GPUs. Our algorithm involves solving a sequence of Ordinary linear Least Squares regression (OLS) problems. For a precise mathematical statement, we recall the seemingly abstract but very convenient definition from [11]; explicit algorithms for the computation of OLS solutions are exposed in §4.1.

Definition 3.1 (Ordinary linear least-squares regression). For $l, l^{\prime} \geq 1$ and for probability spaces $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$ and $\left(\mathbb{R}^{l}, \mathcal{B}\left(\mathbb{R}^{l}\right), \eta\right)$, let $S$ be a $\tilde{\mathcal{F}} \otimes \mathcal{B}\left(\mathbb{R}^{l}\right)$-measurable $\mathbb{R}^{l^{\prime}}$-valued function such that $S(\omega, \cdot) \in \mathbf{L}_{2}\left(\mathcal{B}\left(\mathbb{R}^{l}\right), \eta\right)$ for $\tilde{\mathbb{P}}$-a.e. $\omega \in \tilde{\Omega}$, and $\mathcal{L}$ a linear vector subspace of $\mathbf{L}_{2}\left(\mathcal{B}\left(\mathbb{R}^{l}\right), \eta\right)$ spanned by deterministic $\mathbb{R}^{l^{\prime}}$-valued functions $\left\{p_{k}(),. k \geq 1\right\}$. The least squares approximation of $S$ in the space $\mathcal{L}$ with respect to $\eta$ is the $\left(\tilde{\mathbb{P}} \times \eta\right.$-a.e.) unique, $\tilde{\mathcal{F}} \otimes \mathcal{B}\left(\mathbb{R}^{l}\right)$-measurable function $S^{\star}$ given by

$$
S^{\star}(\omega, \cdot)=\arg \inf _{\phi \in \mathcal{L}} \int|\phi(x)-S(\omega, x)|^{2} \eta(\mathrm{~d} x)
$$

We say that $S^{\star}$ solves $\mathbf{O L S}(S, \mathcal{L}, \eta)$.
On the other hand, suppose that $\eta_{M}=\frac{1}{M} \sum_{m=1}^{M} \delta_{\mathcal{X}(m)}$ is a discrete probability measure on $\left(\mathbb{R}^{l}, \mathcal{B}\left(\mathbb{R}^{l}\right)\right)$, where $\delta_{x}$ is the Dirac measure on $x$ and $\mathcal{X}^{(1)}, \ldots, \mathcal{X}^{(M)}$ : $\tilde{\Omega} \rightarrow \mathbb{R}^{l}$ are i.i.d. random variables. For an $\tilde{\mathcal{F}} \otimes \mathcal{B}\left(\mathbb{R}^{l}\right)$-measurable $\mathbb{R}^{l^{\prime}}$-valued function $S$ such that $\left|S\left(\omega, \mathcal{X}^{(m)}(\omega)\right)\right|<\infty$ for any $m$ and $\tilde{\mathbb{P}}$-a.e. $\omega \in \tilde{\Omega}$, the least squares approximation of $S$ in the space $\mathcal{L}$ with respect to $\eta_{M}$ is the $\left(\tilde{\mathbb{P}}\right.$-a.e.) unique, $\tilde{\mathcal{F}} \otimes \mathcal{B}\left(\mathbb{R}^{l}\right)$ measurable function $S^{\star}$ given by

$$
S^{\star}(\omega, \cdot)=\arg \inf _{\phi \in \mathcal{L}} \frac{1}{M} \sum_{m=1}^{M}\left|\phi\left(\mathcal{X}^{(m)}(\omega)\right)-S\left(\omega, \mathcal{X}^{(m)}(\omega)\right)\right|^{2}
$$

We say that $S^{\star}$ solves $\mathbf{O L S}\left(S, \mathcal{L}, \eta_{M}\right)$.
Definition 3.2 (Simulations and empirical measures). Recall the Markov chain $\left(X_{j}^{\left(i, \nu_{k}\right)}\right)_{i \leq j \leq N}$ initialized as in $\left(\mathbf{A}_{\nu}\right)$. For any $i \in\{0, \ldots, N-1\}$ and $k \in\{1, \ldots, K\}$, define $M \geq \operatorname{dim}\left(\mathcal{L}_{Y, k}\right) \vee \operatorname{dim}\left(\mathcal{L}_{Z, k}\right)$ independent copies of $\left(\Delta W_{i},\left(X_{j}^{i, \nu_{k}}\right)_{i \leq j \leq N}\right)$ that we denote by

$$
\mathcal{C}_{i, k}:=\left\{\left(\Delta W_{i}^{(i, k, m)},\left(X_{j}^{(i, k, m)}\right)_{i \leq j \leq N}\right): m=1, \ldots, M\right\}
$$

The random variables $\mathcal{C}_{i, k}$ form a cloud of simulations used for the regression at time $i$ and in the stratum $k$. Furthermore, we assume that the clouds of simulations $\left(\mathcal{C}_{i, k}: 0 \leq i \leq N-1,1 \leq k \leq K\right)$ are independently generated. All these random variables are defined on a probability space $\left(\Omega^{(M)}, \mathcal{F}^{(M)}, \mathbb{P}^{(M)}\right)$. Denote by $\nu_{i, k, M}$ the empirical probability measure of the $\mathcal{C}_{i, k}$-simulations, i.e.

$$
\nu_{i, k, M}=\frac{1}{M} \sum_{m=1}^{M} \delta_{\left(\Delta W_{i}^{(i, k, m)}, X_{i}^{(i, k, m)}, \ldots, X_{N}^{(i, k, m)}\right)}
$$

Denoting by $(\Omega, \mathcal{F}, \mathbb{P})$ the probability space supporting $\left(\Delta W_{i}, X^{i, \nu_{k}}: 0 \leq i \leq N-1,1 \leq\right.$ $k \leq K)$, which serves as a generic element for the clouds of simulations $\mathcal{C}_{i, k}$, the full probability space used to analyze our algorithm is the product space $(\bar{\Omega}, \overline{\mathcal{F}}, \overline{\mathbb{P}})=$ $(\Omega, \mathcal{F}, \mathbb{P}) \otimes\left(\Omega^{(M)}, \mathcal{F}^{(M)}, \mathbb{P}^{(M)}\right)$. By a slight abuse of notation, we write $\mathbb{P}$ (resp. $\mathbb{E}$ ) to mean $\overline{\mathbb{P}}$ (resp. $\overline{\mathbb{E}}$ ) from now on.

We now come to the definition of the stratified LSMDP algorithm, which computes random approximations $y_{i}^{(M)}($.$) and z_{i}^{(M)}($.

Algorithm 4 (SRMDP). Recall the linear spaces $\mathcal{L}_{Y, k}$ and $\mathcal{L}_{Z, k}$ from ( $\mathbf{A}_{\text {Strat. }}$ ), the bounds (2.3) and the truncation function $\mathcal{T}_{L}$ (see (1.4)).
Initialization. Set $y_{N}^{(M)}(\cdot):=g(\cdot)$.
Backward iteration for $i=N-1$ to $i=0$. For any stratum index $k \in\{1, \ldots, K\}$, generate the empirical measure $\nu_{i, k, M}$ as in Definition 3.2, and define

$$
\begin{cases}\psi_{Z, i, k}^{(M)}(\cdot) & \text { solution of } \quad \mathbf{O L S}\left(S_{Z, i}^{(M)}\left(w, \underline{\mathbf{x}}_{i}\right), \mathcal{L}_{Z, k}, \nu_{i, k, M}\right)  \tag{3.1}\\ \text { for } & S_{Z, i}^{(M)}\left(w, \underline{\mathbf{x}}_{i}\right):=\frac{1}{\Delta_{t}} S_{Y, i+1}^{(M)}\left(\underline{\mathbf{x}}_{i}\right) w, \\ \left.z_{i}^{(M)}(\cdot)\right|_{\mathcal{H}_{k}}:=\mathcal{T}_{C_{z}}\left(\psi_{Z, i, k}^{(M)}(\cdot)\right) \quad \text { (truncation), } \\ \psi_{Y, i, k}^{(M)}(\cdot) & \text { solution of } \quad \mathbf{O L S}\left(S_{Y, i}^{(M)}\left(\underline{\mathbf{x}}_{i}\right), \mathcal{L}_{Y, k}, \nu_{i, k, M}\right) \\ \text { for } & S_{Y, i}^{(M)}\left(\underline{\mathbf{x}}_{i}\right):=g\left(x_{N}\right)+\sum_{j=i}^{N-1} f_{j}\left(x_{j}, y_{j+1}^{(M)}\left(x_{j+1}\right), z_{j}^{(M)}\left(x_{j}\right)\right) \Delta_{t} \\ \left.y_{i}^{(M)}(\cdot)\right|_{\mathcal{H}_{k}}:=\mathcal{T}_{C_{y}}\left(\psi_{Y, i, k}^{(M)}(\cdot)\right) \quad \text { (truncation) }\end{cases}
$$

where $w \in \mathbb{R}^{q}$ and $\underline{\mathbf{x}}_{i}=\left(x_{i}, \ldots, x_{N}\right) \in\left(\mathbb{R}^{d}\right)^{N-i+1}$.
An important difference between SRMDP and established Monte Carlo algorithms $[7,15,10,11]$ is that the number of simulations falling in each hypercube is no more random but fixed and equal to $M$. Observe first that this is likely to improve the numerical stability of the regression algorithm: there is no risk that too few simulations will land in the hypercube, leading to under-fitting. Later, in $\S 4$, we shall explain how to implement Algorithm 4 on a GPU device. The key point is that the calculations at
every time point are fully independent between the different hypercubes, so that we can perform them in parallel across the hypercubes. The choice of $M$ independent on $k$ is made in order to maintain a computational effort equal on each of the strata. In this way, the gain in parallelization is likely to be the largest. However, the subsequent mathematical analysis can be easily adapted to make the number of simulations vary with $k$ whenever necessary.

An easy but important consequence of Algorithm 4 and of the bounds of Proposition 2.2 is the following absolute bound; the proof is analogous to that of [11, Lemma 4.7].

Lemma 3.3. With the above notation, we have

$$
\sup _{0 \leq i \leq N} \sup _{\underline{\mathbf{x}_{\mathbf{i}}} \in\left(\mathbb{R}^{d}\right)^{N-i+1}}\left|S_{Y, i}^{(M)}\left(\underline{\mathbf{x}_{\mathbf{i}}}\right)\right| \leq C_{3.3}:=C_{g}+T\left(L_{f} C_{y}\left[1+\frac{\sqrt{q}}{\sqrt{\Delta_{t}}}\right]+C_{f}\right)
$$

3.2. Error analysis. The analysis will be performed according to several $\mathbf{L}_{2^{-}}$ norms, either w.r.t. the probability measure $\nu$, or the empirical norm related to the cloud simulations. They are defined as follows:

$$
\begin{aligned}
\mathcal{E}(Y, M, i) & :=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) \mathbb{E}\left(\left|y_{i}^{(M)}(\cdot)-y_{i}(\cdot)\right|_{i, k, M}^{2}\right) \\
\overline{\mathcal{E}}(Y, M, i) & :=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) \mathbb{E}\left(\left|y_{i}^{(M)}(\cdot)-y_{i}(\cdot)\right|_{\nu_{k}}^{2}\right)=\mathbb{E}\left(\left|y_{i}^{(M)}(\cdot)-y_{i}(\cdot)\right|_{\nu}^{2}\right) \\
\mathcal{E}(Z, M, i) & :=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) \mathbb{E}\left(\left|z_{i}^{(M)}(\cdot)-z_{i}(\cdot)\right|_{i, k, M}^{2}\right) \\
\overline{\mathcal{E}}(Z, M, i) & :=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) \mathbb{E}\left(\left|z_{i}^{(M)}(\cdot)-z_{i}(\cdot)\right|_{\nu_{k}}^{2}\right)=\mathbb{E}\left(\left|z_{i}^{(M)}(\cdot)-z_{i}(\cdot)\right|_{\nu}^{2}\right)
\end{aligned}
$$

where

$$
|h|_{i, k, M}:=\left(\int|h|^{2}\left(\omega, \underline{\mathbf{x}}_{i}\right) \nu_{i, k, M}\left(\mathrm{~d} \omega, \mathrm{~d} \underline{\mathbf{x}}_{i}\right)\right)^{1 / 2}
$$

In fact, the norms $\mathcal{E}(., M, i)$ and $\overline{\mathcal{E}}(., M, i)$ are related through model-free concentration-of-measures inequalities. This relation is summarized in the proposition below.

Proposition 3.4. For each $i \in\{0, \ldots, N-1\}$, we have

$$
\begin{aligned}
& \overline{\mathcal{E}}(Y, M, i) \leq 2 \mathcal{E}(Y, M, i)+\frac{2028 C_{y}^{2} \log (3 M)}{M}\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)+1\right) \\
& \overline{\mathcal{E}}(Z, M, i) \leq 2 \mathcal{E}(Z, M, i)+\frac{2028 q C_{y}^{2} \log (3 M)}{\Delta_{t} M}\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right)
\end{aligned}
$$

Proof. It is clearly sufficient to show that

$$
\begin{aligned}
\mathbb{E}\left(\left|y_{i}^{(M)}(\cdot)-y_{i}(\cdot)\right|_{\nu_{k}}^{2}\right) \leq & 2 \mathbb{E}\left(\left|y_{i}^{(M)}(\cdot)-y_{i}(\cdot)\right|_{i, k, M}^{2}\right) \\
& +\frac{2028 C_{y}^{2} \log (3 M)}{M}\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)+1\right)
\end{aligned}
$$

$$
\begin{aligned}
& \mathbb{E}\left(\left|z_{i}^{(M)}(\cdot)-z_{i}(\cdot)\right|_{\nu_{k}}^{2}\right) \leq 2 \mathbb{E} \\
&\left(\left|z_{i}^{(M)}(\cdot)-z_{i}(\cdot)\right|_{i, k, M}^{2}\right) \\
&+\frac{2028 q C_{y}^{2} \log (3 M)}{\Delta_{t} M}\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)+1\right)
\end{aligned}
$$

which follows exactly as in the proof of [11, Proposition 4.10].
From the previous proposition, the controls on $\overline{\mathcal{E}}(Y, M, i)$ and $\overline{\mathcal{E}}(Z, M, i)$ stem from those on $\mathcal{E}(Y, M, i)$ and $\mathcal{E}(Z, M, i)$, which are handled in Theorem 3.5 below.
In order to study the impact of basis selection, we define the squared quadratic approximation errors associated to the basis in hypercube $\mathcal{H}_{k}$ by

$$
T_{i, k}^{Y}:=\inf _{\phi \in \mathcal{L}_{Y, k}}\left|\phi-y_{i}\right|_{\nu_{k}}^{2}, \quad T_{i, k}^{Z}:=\inf _{\phi \in \mathcal{L}_{Z, k}}\left|\phi-z_{i}\right|_{\nu_{k}}^{2} .
$$

These terms are the minimal error that can possibly be achieved by the basis $\mathcal{L}_{Y, k}$ (resp. $\mathcal{L}_{Z, k}$ ) in order to approximate the restriction $\left.y_{i}(\cdot)\right|_{\mathcal{H}_{k}}\left(\right.$ resp. $\left.z_{i}(\cdot)\right|_{\mathcal{H}_{k}}$ ) in the $\mathbf{L}_{2}$ norm. Consequently, the global squared quadratic approximation error is given by

$$
\begin{align*}
T_{i}^{Y} & :=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) T_{i, k}^{Y}=\inf _{\phi \text { s.t. }\left.\phi\right|_{\mathcal{H}_{k}} \in \mathcal{L}_{Y, k}}\left|\phi-y_{i}\right|_{\nu}^{2}  \tag{3.2}\\
T_{i}^{Z} & :=\sum_{k=1}^{K} \nu\left(\mathcal{H}_{k}\right) T_{i, k}^{Z}=\inf _{\phi \text { s.t. } \phi \mid \mathcal{H}_{k} \in \mathcal{L}_{Z, k}}\left|\phi-z_{i}\right|_{\nu}^{2} \tag{3.3}
\end{align*}
$$

As we shall see in Theorem 3.5 below, the terms $T_{i}^{Y}$ and $T_{i}^{Z}$ are closely associated to the limit of the expected quadratic error of the numerical scheme in the asymptotic $M \rightarrow \infty$; for this reason, these terms are usually called bias terms.

Now, we are in the position to state our main result giving non-asymptotic error estimates.

Theorem 3.5 (Error for the Stratified LSMDP scheme). Recall the constants $C_{y}$ from Proposition 2.2, $C_{3.3}$ from Lemma 3.3, and $c_{\left(\mathbf{A}_{\nu}\right)}$ from Proposition 2.1. For each $i \in\{0, \ldots, N-1\}$, define

$$
\begin{aligned}
\mathcal{E}(i):= & 2 \sum_{j=i}^{N-1} \Delta_{t}\left(T_{j}^{Y}+3 C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)}{M}+12168 L_{f}^{2} \Delta_{t} \frac{\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right) q C_{y}^{2} \log (3 M)}{M}\right. \\
& \left.+3 T_{j}^{Z}+6 q C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)}{\Delta_{t} M}\right) \\
& +\left(T-t_{i}\right) \frac{1014 C_{y}^{2} \log (3 M)}{M}\left(\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)+1\right)+\frac{q}{\Delta_{t}}\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right)\right) .
\end{aligned}
$$

For $\Delta_{t}$ small enough such that $L_{f} \Delta_{t} \leq \sqrt{\frac{2}{15}}$ and $\Delta_{t} L_{f}^{2} \leq \frac{1}{288 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}(1+T)}$, we have, for all $0 \leq i \leq N-1$,

$$
\begin{align*}
\mathcal{E}(Y, M, i) \leq & T_{i}^{Y}+3 C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)}{M}+12168 L_{f}^{2} \Delta_{t} \frac{\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right) q C_{y}^{2} \log (3 M)}{M} \\
& +\left(1+15 L_{f}^{2} \Delta_{t}\right) C_{3.5} \mathcal{E}(i), \tag{3.4}
\end{align*}
$$

(3.5) $\sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Z, M, j) \leq C_{3.5} \mathcal{E}(i)$,
where $C_{3.5}:=\exp \left(288 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}(1+T) L_{f}^{2} T\right)$.
3.3. Proof of Theorem 3.5. We start by obtaining estimates on the local empirical quadratic errors terms

$$
\mathbb{E}\left(\left|y_{i}^{(M)}(\cdot)-y_{i}(\cdot)\right|_{i, k, M}^{2}\right), \quad \mathbb{E}\left(\left|z_{i}^{(M)}(\cdot)-z_{i}(\cdot)\right|_{i, k, M}^{2}\right)
$$

on each of the hypercubes $\mathcal{H}_{k}(k=1, \ldots, K)$. We first reformulate (2.2) with $\eta=\nu_{k}$ in terms of the Definition 3.1 of OLS. For each $i \in\{0, \ldots, N-1\}$ and $k \in\{1, \ldots, K\}$, let $\nu_{i, k}:=\mathbb{P} \circ\left(\Delta W_{i}, X_{i}^{i, \nu_{k}}, \ldots, X_{N}^{i, \nu_{k}}\right)^{-1}$, so that we have

$$
\left\{\begin{array}{l}
\left.y_{i}(\cdot)\right|_{\mathcal{H}_{k}} \text { solution of } \mathbf{O L S}\left(S_{Y, i}\left(\underline{\mathbf{x}}_{i}\right), \mathcal{L}_{k}^{(1)}, \nu_{i, k}\right) \\
\text { where } S_{Y, i}\left(\underline{\mathbf{x}}_{i}\right):=g\left(x_{N}\right)+\sum_{j=i}^{N-1} f_{j}\left(x_{j}, y_{j+1}\left(x_{j+1}\right), z_{j}\left(x_{j}\right)\right) \Delta_{t} \\
\left.z_{i}(\cdot)\right|_{\mathcal{H}_{k}} \text { solution of } \mathbf{O L S}\left(S_{Z, i}\left(w, \underline{\mathbf{x}}_{i}\right), \mathcal{L}_{k}^{(q)}, \nu_{i, k}\right) \\
\quad \text { where } S_{Z, i}\left(w, \underline{\mathbf{x}}_{i}\right):=\frac{1}{\Delta_{t}} S_{Y, i+1}\left(\underline{\mathbf{x}}_{i}\right) w,
\end{array}\right.
$$

where $w \in \mathbb{R}^{q}, \underline{\mathbf{x}}_{i}:=\left(x_{i}, \ldots, x_{N}\right) \in\left(\mathbb{R}^{d}\right)^{N-i+1}$ and where $\mathcal{L}_{k}^{\left(l^{\prime}\right)}$ is any dense separable subspace in the $\mathbb{R}^{l^{\prime}}$-valued functions belonging to $\mathbf{L}_{2}\left(\mathcal{B}\left(\mathcal{H}_{k}\right), \nu_{k}\right)$. The above OLS solutions and those defined in (3.1) will be compared with other intermediate OLS solutions given by

$$
\left\{\begin{array}{lll}
\psi_{Y, i, k}(\cdot) & \text { solution of } & \operatorname{OLS}\left(S_{Y, i}\left(\underline{x}_{i}\right), \mathcal{L}_{Y, k}, \nu_{i, k, M}\right) \\
\psi_{Z, i, k}(\cdot) & \text { solution of } & \operatorname{OLS}\left(S_{Z, i}\left(w, \underline{\mathbf{x}}_{i}\right), \mathcal{L}_{Z, k}, \nu_{i, k, M}\right)
\end{array}\right.
$$

In order to handle the dependence on the simulation clouds, we define the following $\sigma$-algebras.

Definition 3.6. Define the $\sigma$-algebras
$\mathcal{F}_{i}^{(*)}:=\sigma\left(\mathcal{C}_{i+1, k}, \ldots, \mathcal{C}_{N-1, k}: 1 \leq k \leq K\right), \quad \mathcal{F}_{i, k}^{(M)}:=\mathcal{F}_{i}^{(*)} \vee \sigma\left(X_{i}^{(i, k, m)}: 1 \leq m \leq M\right)$.
For every $i \in\{0, \ldots, N-1\}$ and $k \in\{1, \ldots, K\}$, let $\mathbb{E}_{i, k}^{(M)}(\cdot)$ (resp. $\left.\mathbb{P}_{i, k}^{M}(\cdot)\right)$ with respect to $\mathcal{F}_{i, k}^{(M)}$. Defining additionally the functions

$$
\begin{aligned}
\xi_{Y, i}^{*}(x) & :=\mathbb{E}\left(S_{Y, i}^{(M)}\left(\underline{\mathbf{X}}_{i}\right)-S_{Y, i}\left(\underline{\mathbf{X}}_{i}\right) \mid X_{i}=x, \mathcal{F}^{(M)}\right) \\
\xi_{Z, i}^{*}(x) & :=\mathbb{E}\left(S_{Z, i}^{(M)}\left(\Delta W_{i}, \underline{\mathbf{X}}_{i}\right)-S_{Z, i}\left(\Delta W_{i}, \underline{\mathbf{X}}_{i}\right) \mid X_{i}=x, \mathcal{F}^{(M)}\right)
\end{aligned}
$$

now we are in the position to prove that

$$
\begin{align*}
\mathbb{E}\left(\left|y_{i}(\cdot)-y_{i}^{(M)}(\cdot)\right|_{i, k, M}^{2}\right) \leq & T_{i, k}^{Y}+6 \mathbb{E}\left(\left|\xi_{Y, i}^{*}(\cdot)\right|_{\nu_{k}}^{2}\right)+3 C_{3.3}^{2} \frac{\operatorname{dim}\left(\mathcal{L}_{Y, k}\right)}{M} \\
& +15 L_{f}^{2} \Delta_{t}^{2} \mathbb{E}\left(\left|z_{i}(\cdot)-z_{i}^{(M)}(\cdot)\right|_{i, k, M}^{2}\right) \\
& +12168 L_{f}^{2} \Delta_{t} \frac{\left(\operatorname{dim}\left(\mathcal{L}_{Z, k}\right)+1\right) q C_{y}^{2} \log (3 M)}{M},  \tag{3.6}\\
\mathbb{E}\left(\left|z_{i}(\cdot)-z_{i}^{(M)}(\cdot)\right|_{i, k, M}^{2}\right) \leq & T_{i, k}^{Z}+2 \mathbb{E}\left(\left|\xi_{Z, i}^{*}(\cdot)\right|_{\nu_{k}}^{2}\right)+2 q C_{3.3}^{2} \frac{\operatorname{dim}\left(\mathcal{L}_{Z, k}\right)}{\Delta_{t} M} . \tag{3.7}
\end{align*}
$$

In fact, the proof of (3.6)-(3.7) follows analogously the proof of $[11,(4.12)-(4.13)]$; in order to follow the steps of that proof, one must note that the term $R_{\pi}$ of that paper is equal to 1 here, $C_{\pi}$ is equal to $\Delta_{t}$, and $\theta_{L}=1$. Moreover, one must exchange all norms, OLS problems, $\sigma$-algebras, and empirical functions from the reference to the localized versions defined in the preceding paragraphs. Indeed, the proof method of [11, (4.12)-(4.13)] is model free in the sense that it does not care about the distribution of the Markov chain at time $t_{i}$.

We now aim at aggregating the previous estimates across the strata and propagating them along time. For this, let

$$
\begin{align*}
& \begin{aligned}
\mathcal{E}_{1}(i):= & \sum_{j=i}^{N-1} \Delta_{t}\left(T_{j}^{Y}+\right.
\end{aligned}+3 C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)}{M}+12168 L_{f}^{2} \Delta_{t} \frac{\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right) q C_{y}^{2} \log (3 M)}{M} \\
&  \tag{3.8}\\
& \\
& \\
&
\end{align*}
$$

where $\Gamma_{i}:=\left(1+\gamma \Delta_{t}\right)^{i}$ with $\gamma$ to be determined below. Next, defining

$$
\begin{equation*}
\gamma:=288 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}(1+T) L_{f}^{2} \tag{3.9}
\end{equation*}
$$

and recalling that $\Delta_{t} L_{f}^{2} \leq \frac{1}{288 c_{\left(\mathbf{A}_{\nu}\right)}^{2}{ }^{C} A .1(1+T)}$, then $\gamma$ and $\Delta_{t}$ satisfy
(3.10) $\max \left(\frac{1}{\gamma} \times 12 c_{\left(\mathbf{A}_{\nu)}\right.}^{2} C_{A .1}(1+T) L_{f}^{2}, \Delta_{t} \times 12 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}(1+T) L_{f}^{2}\right) \leq \frac{1}{6} \times \frac{1}{4}$.

Additionally, $\Gamma_{i} \leq \exp (\gamma T):=C_{3.5}$ for every $0 \leq i \leq N$. Now, multiply (3.6) and (3.7) by $\nu\left(\mathcal{H}_{k}\right) \Delta_{t} \Gamma_{i}$ and sum them up over $i$ and $k$ to ascertain that

$$
\begin{aligned}
& \sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Y, M, j) \Gamma_{j}+\sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Z, M, j) \Gamma_{j} \\
& \quad \leq \sum_{j=i}^{N-1} \Delta_{t}\left(T_{j}^{Y}+3 C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, k}\right)\right)}{M}+12168 L_{f}^{2} \Delta_{t} \frac{\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, k}\right)\right)+1\right) q C_{y}^{2} \log (3 M)}{M}\right) \Gamma_{j} \\
& +\sum_{j=i}^{N-1} \Delta_{t}\left\{\left(T_{j}^{Z}+2 q C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, k}\right)\right)}{\Delta_{t} M}+2 \mathbb{E}\left(\left|\xi_{Z, j}^{*}(\cdot)\right|_{\nu}^{2}\right)\right)\left(1+15 L_{f}^{2} \Delta_{t}^{2}\right)+6 \mathbb{E}\left(\left|\xi_{Y, j}^{*}(\cdot)\right|_{\nu}^{2}\right)\right\} \Gamma_{j} \\
& (3.11) \leq \mathcal{E}_{1}(i)+6 \sum_{j=i}^{N-1} \Delta_{t}\left(\mathbb{E}\left(\left|\xi_{Y, j}^{*}(\cdot)\right|_{\nu}^{2}\right)+\mathbb{E}\left(\left|\xi_{Z, j}^{*}(\cdot)\right|_{\nu}^{2}\right)\right) \Gamma_{j},
\end{aligned}
$$

where we have used $\left(1+15 L_{f}^{2} \Delta_{t}^{2}\right) \leq 3$ (since $L_{f} \Delta_{t} \leq \sqrt{\frac{2}{15}}$ ), and the term $\mathcal{E}_{1}$ from (3.8) above. Next, from Proposition 2.1, we have

$$
\mathbb{E}\left(\left|\xi_{Y, j}^{*}(\cdot)\right|_{\nu}^{2}\right)+\mathbb{E}\left(\left|\xi_{Z, j}^{*}(\cdot)\right|_{\nu}^{2}\right) \leq c_{\left(\mathbf{A}_{\nu}\right)}\left(\mathbb{E}\left(\left|\xi_{Y, j}^{*}\left(X_{j}^{0, \nu}\right)\right|^{2}\right)+\mathbb{E}\left(\left|\xi_{Z, j}^{*}\left(X_{j}^{0, \nu}\right)\right|^{2}\right)\right)
$$

Furthermore, note that $\left(\xi_{Y, j}^{*}\left(X_{j}^{0, \nu}\right), \xi_{Z, j}^{*}\left(X_{j}^{0, \nu}\right): 0 \leq j \leq N-1\right)$ solves a discrete BSDE (in the sense of Appendix A.2) with terminal condition 0 and driver

$$
f_{\xi^{*}, j}(y, z):=f_{j}\left(X_{j}^{0, \nu}, y_{j+1}^{(M)}\left(X_{j+1}^{0, \nu}\right), z_{j}^{(M)}\left(X_{j}^{0, \nu}\right)\right)-f_{j}\left(X_{j}^{0, \nu}, y_{j+1}\left(X_{j+1}^{0, \nu}\right), z_{j}\left(X_{j}^{0, \nu}\right)\right)
$$

This allows the application of Proposition A.1, with the first $\operatorname{BSDE}\left(\xi_{Y, j}^{*}\left(X_{j}^{0, \nu}\right), \xi_{Z, j}^{*}\left(X_{j}^{0, \nu}\right)\right.$ : $0 \leq j \leq N-1$ ), and the second one equal to 0 : since $L_{f_{2}}=0$, any choice of $\gamma>0$ is valid and we take $\gamma$ as in (3.9). We obtain

$$
\begin{aligned}
& \sum_{j=i}^{N-1} \Delta_{t}\left(\mathbb{E}\left(\left|\xi_{Y, j}^{*}(\cdot)\right|_{\nu}^{2}\right)+\mathbb{E}\left(\left|\xi_{Z, j}^{*}(\cdot)\right|_{\nu}^{2}\right)\right) \Gamma_{j} \\
& \leq 6 c_{\left(\mathbf{A}_{\nu}\right)} C_{A .1}(1+T)\left(\frac{1}{\gamma}+\Delta_{t}\right) L_{f}^{2} \sum_{j=i}^{N-1} \Delta_{t} \\
& \times\left[\mathbb{E}\left(\left|y_{j+1}^{(M)}\left(X_{j+1}^{0, \nu}\right)-y_{j+1}\left(X_{j+1}^{0, \nu}\right)\right|^{2}\right)+\mathbb{E}\left(\left|z_{j}^{(M)}\left(X_{j}^{0, \nu}\right)-z_{j}\left(X_{j}^{0, \nu}\right)\right|^{2}\right)\right] \Gamma_{j}
\end{aligned}
$$

Now, Proposition 2.1 yields to

$$
\begin{aligned}
& \mathbb{E}\left(\left|y_{j+1}^{(M)}\left(X_{j+1}^{0, \nu}\right)-y_{j+1}\left(X_{j+1}^{0, \nu}\right)\right|^{2}\right)+\mathbb{E}\left(\left|z_{j}^{(M)}\left(X_{j}^{0, \nu}\right)-z_{j}\left(X_{j}^{0, \nu}\right)\right|^{2}\right) \\
& \leq c_{\left(\mathbf{A}_{\nu}\right)}[\overline{\mathcal{E}}(Y, M, j+1)+\overline{\mathcal{E}}(Z, M, j)] \\
& \leq 2 c_{\left(\mathbf{A}_{\nu}\right)}[\mathcal{E}(Y, M, j+1)+\mathcal{E}(Z, M, j)]+c_{\left(\mathbf{A}_{\nu}\right)} \frac{2028 C_{y}^{2} \log (3 M)}{M}\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)+1\right) \\
& \quad+c_{\left(\mathbf{A}_{\nu}\right)} \frac{2028 q C_{y}^{2} \log (3 M)}{\Delta_{t} M}\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right)
\end{aligned}
$$

where the last inequality follows from the concentration-measure inequalities in Proposition 3.4. In order to summarize this, we define
$\mathcal{E}_{2}(i):=\frac{1014 C_{y}^{2} \log (3 M)}{M}\left(\sum_{j=i}^{N-1} \Delta_{t} \Gamma_{j}\right)\left(\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)+1\right)+\frac{q}{\Delta_{t}}\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right)\right)$
and make use of (3.10), and that $\Gamma_{j} \leq \Gamma_{j+1}$ in order to ascertain that we have

$$
\begin{aligned}
& \sum_{j=i}^{N-1} \Delta_{t}\left(\mathbb{E}\left(\left|\xi_{Y, j}^{*}(\cdot)\right|_{\nu}^{2}\right)+\mathbb{E}\left(\left|\xi_{Z, j}^{*}(\cdot)\right|_{\nu}^{2}\right)\right) \Gamma_{j} \\
& \leq 12 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}(1+T)\left(\frac{1}{\gamma}+\Delta_{t}\right) L_{f}^{2}\left[\sum_{j=i}^{N-1} \Delta_{t}(\mathcal{E}(Y, M, j)+\mathcal{E}(Z, M, j)) \Gamma_{j}+\mathcal{E}_{2}(i)\right] \\
& \leq \frac{1}{6} \times \frac{1}{2}\left[\sum_{j=i}^{N-1} \Delta_{t}(\mathcal{E}(Y, M, j)+\mathcal{E}(Z, M, j)) \Gamma_{j}+\mathcal{E}_{2}(i)\right]
\end{aligned}
$$

By plugging this into (3.11) readily yields to

$$
\begin{aligned}
\sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Y, M, j) \Gamma_{j} & +\sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Z, M, j) \Gamma_{j} \\
& \leq \mathcal{E}_{1}(i)+\frac{1}{2}\left[\sum_{j=i}^{N-1} \Delta_{t}(\mathcal{E}(Y, M, j)+\mathcal{E}(Z, M, j)) \Gamma_{j}+\mathcal{E}_{2}(i)\right]
\end{aligned}
$$

and therefore

$$
\begin{equation*}
\sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Y, M, j) \Gamma_{j}+\sum_{j=i}^{N-1} \Delta_{t} \mathcal{E}(Z, M, j) \Gamma_{j} \leq 2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i) \tag{3.12}
\end{equation*}
$$

This completes the proof of the estimate (3.5) on $z$ as stated in Theorem 3.5, using $1 \leq \Gamma_{i} \leq C_{3.5}$ and $2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i) \leq C_{3.5} \mathcal{E}(i)$. It remains to derive (3.4). Starting from (3.6), multiplying by $\nu\left(\mathcal{H}_{k}\right)$ and summing over $k$ yields to

$$
\begin{align*}
\mathcal{E}(Y, M, i) \leq & T_{i}^{Y}+6 \mathbb{E}\left(\left|\xi_{Y, i}^{*}(\cdot)\right|_{\nu}^{2}\right)+3 C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)}{M} \\
& +15 L_{f}^{2} \Delta_{t}\left(2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i)\right) \\
& +12168 L_{f}^{2} \Delta_{t} \frac{\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right) q C_{y}^{2} \log (3 M)}{M} \tag{3.13}
\end{align*}
$$

where we use the inequality (3.12) to control $\Delta_{t} \mathcal{E}(Z, M, i)$. Using the same arguments as before, we upper bound $\mathbb{E}\left(\left|\xi_{Y, i}^{*}(\cdot)\right|_{\nu}^{2}\right)$ by

$$
6 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}\left(\frac{1}{\gamma}+\Delta_{t}\right) L_{f}^{2} \sum_{j=i}^{N-1} \Delta_{t}(\overline{\mathcal{E}}(Y, M, j)+\overline{\mathcal{E}}(Z, M, j)) \Gamma_{j}
$$

By additionally bounding $\overline{\mathcal{E}}(Y, M, j)$ and $\overline{\mathcal{E}}(Z, M, j)$ using the concentration-measure inequalities of Proposition 3.4 and plugging this in (3.13), we finally obtain

$$
\begin{aligned}
\mathcal{E}(Y, M, i) & \leq T_{1, i}^{Y}+3 C_{3.3}^{2} \frac{\nu\left(\operatorname{dim}\left(\mathcal{L}_{Y, .}\right)\right)}{M}+15 L_{f}^{2} \Delta_{t}\left(2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i)\right) \\
& +12168 L_{f}^{2} \Delta_{t} \frac{\left(\nu\left(\operatorname{dim}\left(\mathcal{L}_{Z, .}\right)\right)+1\right) q C_{y}^{2} \log (3 M)}{M} \\
& +72 c_{\left(\mathbf{A}_{\nu}\right)}^{2} C_{A .1}\left(\frac{1}{\gamma}+\Delta_{t}\right) L_{f}^{2}\left[\sum_{j=i}^{N-1} \Delta_{t}(\mathcal{E}(Y, M, j)+\mathcal{E}(Z, M, j)) \Gamma_{j}+\mathcal{E}_{2}(i)\right] .
\end{aligned}
$$

From (3.10) and (3.12), the last term in previous inequality is bounded by

$$
\left(\frac{1}{4(1+T)}+\frac{1}{4(1+T)}\right)\left(2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i)+\mathcal{E}_{2}(i)\right) \leq \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i) \leq 2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i)
$$

This completes the proof of (3.6), using again $2 \mathcal{E}_{1}(i)+\mathcal{E}_{2}(i) \leq C_{3.5} \mathcal{E}(i)$.
4. GPU implementation. In this section, we consider the computation of $y_{i}^{(M)}(\cdot)$ for a given stratum $\mathcal{H}_{k}$ and time point $i$. The calculation of $z_{i}^{(M)}(\cdot)$ is rather similar, only requiring component-wise calculations to be taken into account, so that we do not provide details. The theoretical description of the calculation was given in $\S 3.1$. In this section, we first describe the required computations to implement the approximations with LP0 and LP1 local polynomials in §4.1, and then present their implementation on the GPU in $\S 4.2$.

### 4.1. Explicit solutions to OLS in Algorithm 4.

LP0. This piecewise solution is given by the simple formula [13, Ch. 4]

$$
\begin{equation*}
\left.y_{i}^{(M)}(\cdot)\right|_{H_{k}}=\mathcal{T}_{C_{y}}\left(\frac{\sum_{m=1}^{M} S_{Y, i}^{(M)}\left(\underline{\mathbf{X}}_{i}^{(i, k, m)}\right)}{M}\right) \tag{4.1}
\end{equation*}
$$

Observe that there will be a local memory consumption of $O(1)$ per hypercube to store the simulations needed for the computation of $S_{Y, i}^{(M)}\left(\underline{\mathbf{X}}_{i}^{(i, k, m)}\right)$. Once added in the sum (4.1), their allocation can be freed.

LP1. Let $A$ be the $\mathbb{R}^{M} \otimes \mathbb{R}^{d+1}$ matrix, the components of which are given by $A[m, j]=1 \mathbf{1}_{\{0\}}(j)+X_{i, j}^{(i, k, m)} \mathbf{1}_{\{0\}^{c}}(j)$, where $X_{i, j}^{(i, k, m)}$ is the $j$-th component of $X_{i}^{(i, k, m)}$, and let $S$ be the $\mathbb{R}^{M}$ vector given by $S[m]=S_{Y, i}^{(M)}\left(\underline{\mathbf{X}}_{i}^{(i, k, m)}\right)$. In order to compute $\left.y_{i}^{(M)}(\cdot)\right|_{H_{k}}$, we first perform a QR-factorization $A=Q R$, where $Q$ is and $\mathbb{R}^{M} \otimes \mathbb{R}^{M}$ orthogonal matrix, and $R$ is an $\mathbb{R}^{M} \otimes \mathbb{R}^{d+1}$ upper triangular matrix. The computational cost to compute this factorization is $(d+1)^{2}(M-(d+1) / 3)$ flops using the Householder reflections method [12, Alg. 5.3.2]. Using the form of LP1 and the density of $\nu_{k}$, we can prove that the rank of $A$ is $d+1$ with probability 1 , i.e. $R$ is invertible a.s. (the OLS problem is non-degenerate).

Then, we obtain the approximation $\left.y_{i}^{(M)}(\cdot)\right|_{\mathcal{H}_{k}}$ by computing the coefficients $\alpha=$ $\left(\alpha_{0}, \ldots, \alpha_{d}\right) \in \mathbb{R}^{d+1}$ using the QR factorization and backward-substitution method as follows:

$$
\begin{equation*}
R \alpha=Q^{\top} S, \quad y_{i}^{(M)}(x(k))=\mathcal{T}_{C_{y}}\left(\alpha_{0}+\sum_{j=1}^{d} \alpha_{j+1} \times x_{j}(k)\right) \tag{4.2}
\end{equation*}
$$

for any vector $\left(x(k)=\left(x_{1}(k), \ldots, x_{d}(k)\right)\right.$ in $H_{k}$. By using the Householder reflection algorithm for computing the QR-factorization, there will be a local memory consumption of $O(M \times(d+1))$ for the storage of the matrix $A$ on each hypercube. This memory can be deallocated once the computation (4.2) is completed. We remark that the memory consumption is considerably lower than other alternative QR-factorization methods, as for example the Givens rotations method [12, Alg. 5.2.2], which requires a memory consumption $O\left(M^{2}\right)$ to store the matrix $Q$. This reduced memory consumption is instrumental in the GPU approach, as we explain in forthcoming §5.2.2.
4.2. Pseudo-algorithms for GPU. Algorithm 4 will be implemented on an NVIDIA GPU device. The device architecture is built around a scalable array of multithreaded Streaming Multiprocessors (SMs); each multiprocessor is designed to execute hundreds of threads concurrently. To manage such a large amount of threads, it employs a unique architecture called SIMT (Single-Instruction, Multiple-Thread). The code execution unit is called a kernel and is executed simultaneously on all SMs by independent blocks of threads. Each thread is assigned to a single processor and executes within its own execution environment. Thus, all threads run the same instruction at a time, although over different data. In this section we briefly describe pseudo-codes for the Algorithm 4.

The algorithm has been programmed using the Compute Unified Device Architecture (CUDA) toolkit, specially designed for NVIDIA GPUs, see [17]. The code was built from an optimized C code. The below pseudo-algorithms reflect this programming feature. For the generation of the random numbers in parallel we took advantage of the NVIDIA CURAND library, see [18].

The time loop corresponding to the backward iteration of Algorithm 4 is shown in Listing 1; the kernel corresponds to the use of either the LP0 or the LP1 basis. In Listing 2, a sketch for the LP0 kernel is given. Notice that we are paralellizing the loop for any stratum index $k \in\{1, \ldots, K\}$ in the Algorithm 4 ; the terms $S_{Y, i}^{(M)}\left(\underline{\mathbf{x}}_{i}\right)$ and $S_{Z, i}^{(M)}\left(w, \underline{\mathbf{x}}_{i}\right)$ are computed in the compute_responses_i function, and the coefficients for $\psi_{Y, i, k}^{(M)}(\cdot)$ and $\psi_{Z, i, k}^{(M)}(\cdot)$ are computed in compute_psi_Y and compute_psi_Z, respectively, according to (4.1). Having in view an optimal performance, memory accesses to the simulations, responses and regression coefficients are coalesced, see [17]. In Listing 3, the sketch for the $\mathbf{L P} 1$ kernel is given. Additionally to the tasks of the kernel in Listing 2, each thread builds the matrix $A$ and applies a QR factorization, as detailed in $\S 4.1$. The coefficients for $\psi_{Y, i, k}^{(M)}(\cdot)$ and $\psi_{Z, i, k}^{(M)}(\cdot)$ are computed according to (4.2).

```
int i
curandState *devStates
Initialize devStates
Initialize n_blocks, n_threads_per_block
for (i=N-1; i > = 0; i --)
    kernel_bsde<<<<n_blocks, n_threads_per_block>>>(i, devStates, ...)
                                    LISTING 1
                                    Backward iteration for i=N-1 to i=0.
```

```
_-global_- void kernel_bsde_LP0(int i, curandState* devStates, ...) {
```

_-global_- void kernel_bsde_LP0(int i, curandState* devStates, ...) {
const unsigned int global_tid = blockDim.x * blockIdx.x + threadIdx.x
const unsigned int global_tid = blockDim.x * blockIdx.x + threadIdx.x
curandState localState = devStates[global_tid]
curandState localState = devStates[global_tid]
unsigned long long int bin
unsigned long long int bin
for(bin=global_tid; bin <K; bin+=n_blocks*n_threads_per_block) {
for(bin=global_tid; bin <K; bin+=n_blocks*n_threads_per_block) {
simulates_x(\&localState, global_tid, bin, ...)
simulates_x(\&localState, global_tid, bin, ...)
compute_responses_i(\&localState, global_tid, i, ...)
compute_responses_i(\&localState, global_tid, i, ...)
compute_psi_Z(global_tid, bin, i, ...)
compute_psi_Z(global_tid, bin, i, ...)
compute_psi_Y(global_tid, bin, i, ...)
compute_psi_Y(global_tid, bin, i, ...)
}
}
devStates[global_tid] = localState
devStates[global_tid] = localState
}

```

LISTING 2
Kernel for the approximation with \(\mathbf{L P 0}\).
```

--global_- void kernel_bsde_LP1(int i, curandState *devStates, ...) {
const unsigned int global_tid = blockDim.x * blockIdx.x + threadIdx.x
curandState localState = devStates[global_tid]
unsigned long long int bin
for(bin=global_tid;bin<K;bin+=n_blocks*n_threads_per_block) {
simulates_x(\&localState, global_tid, bin, ...)
compute_responses_i(\&localState, global_tid, i, ...)
build_d_A(global_tid, d_A, ...)
qr(global_tid, d_A, ...)
compute_psi_Z(global_tid, bin, i, d_A, ...)
compute_psi_Y(global_tid, bin, i, d_A, ...)
}
devStates[global_tid] = localState
}

```
4.3. Theoretical complexity analysis. In this section, we assume that the functions \(y_{i}(\cdot)\) and \(z_{i}(\cdot)\) are smooth, namely globally Lipschitz (resp. \(C^{1}\) and the first derivatives are globally Lipschitz) in the LP0 (resp. LP1) case. The strata will be composed of uniform hypercubes of side length \(\delta>0\) in the domain \([-L, L]^{d}\), where \(L=\log (N) / \mu\) and \(\mu\) is the parameter of the logistic distribution. This choice ensures \(\nu\left(\mathbb{R}^{d} \backslash[-L, L]^{d}\right) \leq 2 d \exp (-\mu L)=O\left(N^{-1}\right)\). Our aim is to calibrate the numerical parameters (number of simulations and number of strata) so that the error given in Theorem 3.5 is \(O\left(N^{-1}\right)\), where \(N\) is the number of time-steps. This tolerance error is the one we usually obtain after time discretization with \(N\) time points [20, 8, 19]. In the following, we focus on polynomial dependency w.r.t. \(N\), keeping only the highest degree, ignoring constants and \(\log (N)\) terms.

Squared bias errors \(T_{1, i}^{Y}\) and \(T_{1, i}^{Z}\) in (3.2)-(3.3). First, we remark that the approximation error of the numerical scheme, namely the error due to basis selection, depends principally on the size \(\delta\) of strata. In the case of \(\mathbf{L P} \mathbf{0}\), the squared bias error is proportional to the squared hypercube diameter plus the tail contribution, i.e. \(O\left(\delta^{2}+\nu\left(\mathbb{R}^{d} \backslash[-L, L]^{d}\right)\right)\); to calibrate this bias to \(O\left(N^{-1}\right)\), we require \(\delta=O\left(N^{-1 / 2}\right)\). In contrast, the squared bias in \([-L, L]^{d}\) using LP1 is proportional to the fourth power of the hypercube diameter, whence \(\delta=O\left(N^{-1 / 4}\right)\). As a result, ignoring the log terms the number of required hypercubes is
\[
\mathbf{L P 0}: \quad K=O\left(N^{d / 2}\right), \quad \mathbf{L P 1}: \quad K=O\left(N^{d / 4}\right)
\]
in both cases.
Statistical and interdependence errors. These error terms depend on the number of local polynomials, as well as on the number of simulations. Indeed, denoting \(K^{\prime}=\operatorname{dim}\left(\mathcal{L}_{Y}\right.\) or \(\left.Z,.\right)\) the number of local polynomials and \(M\) the number of simulations in the hypercube, then both errors are dominated by \(O\left(N K^{\prime} \log (M) / M\right)\), where the factor \(N\) comes from the \(Z\) part of the solution (see \(\mathcal{E}(i)\) in Theorem 3.5). For LP0 (resp. LP1), \(K^{\prime}=1\) or \(q\) (resp. \(K^{\prime}=d+1\) or \(q(d+1)\) ). This implies to select
\[
\text { LP0 : } \quad M=O\left(N^{2}\right), \quad \text { LP1 : } \quad M=O\left(N^{2}\right)
\]
again omitting the log terms.
Computational cost. The computational cost (in flops) of the simulations per hypercube is equal to \(O(M \times N)\), because we simulate \(M\) paths (of length \(N\) ) of the process \(X\). The cost of the regression per hypercube is \(O(M \times N)\), see \(\S 4.1\), and thus equivalent to the simulation cost. Putting in the values of \(M\) from the last paragraph, the overall computational cost \(\mathcal{C}_{\text {cost }}\) (summed over all hypercubes and time steps) is
\[
\mathbf{L P 0}: \quad \mathcal{C}_{\mathrm{cost}}^{\mathrm{CPU}}=O\left(N^{4+d / 2}\right), \quad \mathbf{L P 1}: \quad \mathcal{C}_{\mathrm{cost}}^{\mathrm{cPU}}=O\left(N^{4+d / 4}\right)
\]

This quantity is related to the computational time for a sequential system (CPU implementation) where there is no parallel computing. For the GPU implementation, described in \(\S 4.2\), there is an additional computational time improvement since the computations on the hypercubes will be threaded across the cores of the card. Thus, the computational cost on GPU is
LP0 : \(\quad \mathcal{C}_{\text {cost }}^{\text {GPU }}=O\left(N^{4+d / 2}\right) / \mathcal{C}_{\text {Load factor }}, \quad \mathbf{L P 1}: \quad \mathcal{C}_{\text {cost }}^{\mathrm{GPU}}=O\left(N^{4+d / 4}\right) / \mathcal{C}_{\text {Load factor }}\).
where the load factor \(\mathcal{C}_{\text {Load factor }}\) is ideally the number of threads on the device.
Finally, we quantify the improvement in memory consumption offered by the SRMDP algorithm compared to the LSMDP algorithm of [11]. This is a very important
improvement, because, as explained in the introduction, the memory is the key constraint in solving problems in high dimension. We only compare sequential versions of the algorithms, meaning that the computational costs will be the same. The main difference between the two schemes is then in the number of simulations that must be stored in the machine at any given time. We summarize this in Table 1 below.
\begin{tabular}{|c|c|c|c|c|}
\hline Algorithm & \multicolumn{2}{|c|}{ Number of } & \multicolumn{2}{c|}{ Computational } \\
& simulations & \multicolumn{2}{c|}{ cost } \\
& LP0 & LP1 & LP0 & LP1 \\
\hline SRMDP & \(N^{2}\) & \(N^{2}\) & \(N^{4+d / 2}\) & \(N^{4+d / 4}\) \\
LSMDP & \(N^{2+d / 2}\) & \(N^{2+d / 4}\) & \(N^{4+d / 2}\) & \(N^{4+d / 4}\) \\
\hline \multicolumn{4}{|c|}{ TABLE 1}
\end{tabular}

Comparison of numerical parameters with or without stratified sampling, as a function of \(N\).

In SRMDP, the shared memory is mainly related to storing coefficients representing the solutions on hypercubes, that is \(O\left(N \times \operatorname{dim}\left(\mathcal{L}_{Y}\right.\right.\) or \(\left.\left.Z,.\right) \times K\right)\); if one is using the LP1 basis, one must also take into account the local memory consumption \(M \times(d+1)=O\left(N^{2}\right)\) for the QR factorization, explained in \(\S 4.1\). In contrast, the memory consumption for LSMDP is mainly \(O\left(K \times N^{2}\right)\), which represents the number of simulated paths of the Markov chains that must be stored in the machine at any given time. We summarize the memory consumption of the two algorithms in Table 2.
\begin{tabular}{|c|c|c|}
\hline Algorithm & LP0 & LP1 \\
\hline SRMDP & \(N^{1+d / 2}\) & \(N^{1+d / 4} \vee N^{2}\) \\
LSMDP & \(N^{2+d / 2}\) & \(N^{2+d / 4}\) \\
\hline \multicolumn{3}{|c|}{ TABLE 2} \\
\hline
\end{tabular}

Comparison of shared memory requirement as a function of \(N\).

Observe that SRMDP requires \(N\) times less memory than LSMDP with the LP0 basis. This implicitly implies a gain of 2 on the dimension \(d\) that can be handled. On the other hand, if the LP1 basis is used, the SRMDP requires \(O\left(N^{d / 4}\right)\) less memory for \(d \leq 4\) than LSMDP, and \(N\) times less memory for \(d \geq 4\). Therefore, there is an implicit gain of 4 in the dimension that can be handled by the algorithm.

\section*{5. Numerical experiments.}
5.1. Model, stratification, and performance benchmark. We use the Brownian motion model \(X=W(d=q)\). Moreover, the numerical experiments will consider the performance according to the dimension \(d\). We introduce the function \(\omega(t, x)=\exp \left(t+\sum_{k=1}^{q} x_{k}\right)\). We perform numerical experiments on the BSDE with data \(g(x)=\omega(T, x)(1+\omega(T, x))^{-1}\) and
\[
f(t, x, y, z)=\left(\sum_{k=1}^{q} z_{k}\right)\left(y-\frac{2+q}{2 q}\right)
\]
where \(z=\left(z_{1}, \ldots, z_{q}\right)\). The BSDE has explicit solutions in this framework, given by
\[
y_{i}(x)=\omega\left(t_{i}, x\right)\left(1+\omega\left(t_{i}, x\right)\right)^{-1}, \quad z_{k, i}(x)=\omega\left(t_{i}, x\right)\left(1+\omega\left(t_{i}, x\right)\right)^{-2}
\]
where \(z_{k, i}(x)\) is the \(k\)-th component of the \(q\)-dimensional cylindrical function \(z_{i}(x) \in\) \(\mathbb{R}^{q}\).

The logistic distribution for Algorithm 4 is parameterized by \(\mu=1\). For the leastsquares Monte Carlo, we stratify the domain \([-6.5,6.5]^{q}\) with uniform hypercubes. To assess the performance of the algorithm, we compute the average mean squared error (MSE) over 10 independent runs of the algorithm for three error indicators:
\[
\begin{aligned}
& M S E_{Y, \text { max }}:=\ln \left\{10^{-3} \max _{0 \leq i \leq N-1} \sum_{m=1}^{10^{3}}\left|y_{i}\left(R_{i, m}\right)-y_{i}^{(M)}\left(R_{i, m}\right)\right|^{2}\right\}, \\
& M S E_{Y, \mathrm{av}}:=\ln \left\{10^{-3} N^{-1} \sum_{m=1}^{10^{3}} \sum_{i=0}^{N-1}\left|y_{i}\left(R_{i, m}\right)-y_{i}^{(M)}\left(R_{i, m}\right)\right|^{2}\right\}, \\
& M S E_{Z, \mathrm{av}}:=\ln \left\{10^{-3} N^{-1} \sum_{m=1}^{10^{3}} \sum_{i=0}^{N-1}\left|z_{i}\left(R_{i, m}\right)-z_{i}^{(M)}\left(R_{i, m}\right)\right|^{2}\right\},
\end{aligned}
\]
where the simulations \(\left\{R_{i, m} ; i=0, \ldots, N-1, m=1, \ldots, 10^{3}\right\}\) are independent and identically \(\nu\)-distributed, and independently drawn from the simulations used for the LSMC scheme. We parameterize the hypercubes according to the instructions given in the theoretical complexity analysis, see \(\S 4.3\). In particular, we consider different values of \(N\) and always set \(K=O\left(N^{d / 2}\right)\) in LP0 (resp. \(O\left(N^{d / 4}\right)\) in LP1) and \(M=O\left(N^{2}\right)\). Note, however, that we do not specify the value of \(\delta\), but rather the number of hypercubes per dimension \(K^{1 / q}\), which we denote \#C in what follows; this being equivalent to setting \(\delta\), but is more convenient to program. As we shall illustrate, the error converges as predicted as \(N\) increases, although the exact error values will depend on the constants that we choose in the parameterization of \(K\) and \(M\).
5.2. CPU and GPU performance. In this section, several experiments based on \(\S 5.1\) are presented to assess the performance of CUDA implementation of Algorithm 4 ; the pseudo-algorithms are given in \(\S 4.2\). We shall compare its performance to a version of SRMDP implemented only on CPU, so that the computations across the hypercubes are performed sequentially. Moreover, in order to test the theoretical results of \(\S 4.3\), we compare the performance of the two algorithms according to the choice of the basis functions, the impact of this choice on the convergence of the approximation of the BSDE , and the impact of this choice on the computational performance in terms of computational time and memory consumption.

There are two types of basis functions we investigate: LP0 in §5.2.1, and LP1 in \(\S 5.2 .2\). As explained in \(\S 4.3\), the \(\mathbf{L P 0}\) basis is highly suited to GPU implementation because it has a very low memory requirement per thread of computation. On the other hand, it has a very high global memory requirement for storing coefficients. This represents a problem in high dimensions because one needs many coefficients to obtain a good accuracy. On the other hand, the LP1 basis involves a higher cost per thread, although requires a far lower global memory for storing coefficients; this implies that the impact of the GPU implementation is lower in moderate dimensional problems, but that one can solve problems in higher dimension. Moreover, the full performance impact of the GPU implementation on the LP1 basis is in high dimension, where the number of strata is very high and therefore the GPU is better saturated with computations. We illustrate numerically all of these effects in the following sections.

The numerical experiments have been performed with the following hardware and software configurations: a GPU GeForce GTX TITAN Black with 6 GBytes of global memory (see [16] for details in the architecture), a recent Intel Xeon CPU E5-2620 v2
clocked at 2.10 GHz with 62 GBytes of RAM, CentOS Linux, NVIDIA CUDA SDK 7.0 and GNU C compiler 4.8.2. The CPU programs were compiled with and without -03 flag (which takes advantage of CPU vectorial extensions); the compilation without -03 flag is closer to the pure sequential version. Since the CUDA code has been derived from an optimized C code, both codes perform exactly the same operations, and their performance can be fairly compared according to computation time alone; the CPU time (CPU), the CPU time using -03 flag (VE) and the GPU time (GPU) will all be measured in seconds in the tables of results tables.
5.2.1. Examples with the approximation with LP0 local polynomials. All examples will be run using a \(256 \times 64\) threads configuration. In Table 3 we show results for \(d=4\), with \(\# C=\lfloor 4 \sqrt{N}\rfloor\) and \(M=N^{2}\). Except for the case \(\Delta_{t}=0.2\) where there are not enough strata to fully take advantage of the GPU, the GPU implementation provides a significant reduction in the computational time: the GPU speed-up reaches the value 521 (without -03 flag). Moreover, the speed-up improves as we increase the \#C. Those cells filled with the symbol \(*\) represent CPU times that are prohibitive in terms of computational time.
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#C & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 8 & 25 & -3.712973 & -3.774071 & -0.964842 & 1.74 & 1.25 & 2.00 \\
\hline 0.1 & 12 & 100 & -4.066741 & -4.303750 & -1.607104 & 112.64 & 78.76 & 2.20 \\
\hline 0.05 & 17 & 400 & -4.337988 & -4.698645 & -2.302092 & 6462.19 & 4352.33 & 12.39 \\
\hline 0.02 & 28 & 2500 & -4.472564 & -4.988069 & -3.225411 & \(*\) & \(*\) & 3070.92 \\
\hline \multicolumn{7}{c|}{ TABLE 3} \\
LPO local polynomials, \(d=4, \# C=\lfloor 4 \sqrt{N}\rfloor, M=N^{2}\).
\end{tabular}

Tables 4 and 5 show results for \(d=6\) with \(\# \mathrm{C}=\lfloor\sqrt{N}\rfloor\) and \(\# \mathrm{C}=\lfloor 2 \sqrt{N}\rfloor\), respectively. Convergence is clearly improved by doubling \#C. In Table 5 the case of \(\Delta_{t}=0.02\) is not shown due to insufficient GPU global memory. In Table 4 , the GPU speed-up reaches 531, whereas in Table 5 it reaches 558 (without -03 compiler flag). As in Table 3, the increase in the speed-up is explained due to the increased number of hypercubes, thus demonstrating how important it is to have many hypercubes in the GPU implementation. However, the finer basis requires \(2^{6}\) times as much memory for storing coefficients.
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#C & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 25 & -2.392320 & -2.451332 & -0.431059 & 0.04 & 0.03 & 1.99 \\
\hline 0.1 & 3 & 100 & -2.440274 & -2.500775 & -1.096603 & 5.77 & 4.19 & 2.05 \\
\hline 0.05 & 4 & 400 & -2.829757 & -2.905192 & -1.687142 & 451.18 & 329.40 & 3.15 \\
\hline 0.02 & 7 & 2500 & -3.235130 & -3.539011 & -2.557686 & 464444.02 & 317954.67 & 874.25 \\
\hline \multicolumn{8}{c|}{ TABLE 4} \\
LP0 local polynomials, \(d=6, \# C=\lfloor\sqrt{N}\rfloor, M=N^{2}\).
\end{tabular}

Table 6 shows that the algorithm can work for \(d=11\) in several seconds with a reasonable accuracy in a GPU. The corresponding speed-up with respect to CPU version is around 400.
5.2.2. Examples with the approximation with LP1 local polynomials. In this section we show the results corresponding to the approximation with the LP1 basis. Compared to LP0, this basis consumes much less global memory to store
\begin{tabular}{|c|c|c|c|c|c|c|c|c|}
\hline \(\Delta_{t}\) & \#C & M & \(M S E_{Y, \text { max }}\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline 0.2 & 4 & 25 & -2.707882 & -2.784022 & -0.477751 & 2.52 & 1.97 & 1.94 \\
\hline 0.1 & 6 & 100 & -3.195937 & -3.294488 & -1.133834 & 374.19 & 263.91 & 2.44 \\
\hline 0.05 & 8 & 400 & -3.505867 & -3.664396 & -1.795697 & 29172.89 & 20998.77 & 52.20 \\
\hline \multicolumn{9}{|c|}{\begin{tabular}{l}
TABLE 5 \\
ls, \(d=6, \# C=\)
\end{tabular}} \\
\hline
\end{tabular}
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#C & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 25 & -2.152253 & -2.202357 & 0.211590 & 2.26 & 1.69 & 1.992 \\
\hline 0.1 & 3 & 100 & -2.144843 & -2.267742 & -0.469759 & 2541.45 & 1908.24 & 6.29 \\
\hline 0.05 & 4 & 400 & -2.484169 & -2.633602 & -1.070096 & \(*\) & \(*\) & 2108.64 \\
\hline \multicolumn{7}{|c|}{ TABLE 6} \\
LPO local polynomials, \(d=11, \# C=\lfloor\sqrt{N}\rfloor, M=N^{2}\).
\end{tabular}
coefficients, because it requires far fewer hypercubes, see \(\S 4.3\). On the other hand, the approximation with LP1 basis demands higher local memory due to the storage of a large matrix for each hypercube, as explained in \(\S 4.3\). This may have an impact on the computational time on the GPU: recalling from \(\S 4.2\) the GPU handles multiple hypercubes at any given time, each one requiring the storage of a matrix \(A\), the global memory capacity of the GPU device restricts the number of threads we can handle at any given time. This issue is much less significant with the LP0 basis. In order to optimize the performance of the LP1 basis, we must minimize the local memory storage. We implement the Householder reflection method for QR-factorization, [12, Alg. 5.3.2]. For this, we must store a matrix containing \(M \times(d+1)=O\left(N^{2}\right)\) floating point values per thread on the GPU memory. Nonetheless, in all subsequent examples, we are able to use the highly efficient \(256 \times 64\) thread configuration on the GPU. This leads to high speed-ups not only compared to the sequential code, but also compared with the LP0 basis. Moreover, thanks to the reduced global memory storage for coefficients, we are able to work in a rather high dimension \(d=16\) whilst maintaining an optimal thread configuration; this was not the case in LP0.

REmARK 5.1. There are many methods to implement QR-factorization. However, the choice of method has a substantial impact on the performance of the GPU implementation. For example, the Givens rotation method [12, Alg. 5.2.2] requires the storage of an \(M \times M\) matrix, which corresponds to \(O\left(N^{4}\right)\) floating points. This is rather more than the required \(O\left(N^{2}\right)\) for the Householder reflection method given in Section 4.1. Therefore, the Givens rotation method would be far slower when implemented on a GPU than the Householder reflection method, because it may not be possible to use an optimal thread configuration.

REMARK 5.2. In the forthcoming examples, we use more simulations per stratum for the LP1 basis compared to the equivalent results for \(\mathbf{L P} \mathbf{0}\). This is to account for the additional statistical and interdependence errors, as explained in §4.3.

In Table 7, we present results for \(d=4\). These results are to be compared with Table 3, where in particular the \(M S E_{Z, \text { av }}\) results are closer line to line. The computational time is substantially improved for the CPU and VE calculations. On the other hand, the GPU performance is better only for \(\Delta_{t}=0.02\). Also note that, unlike for the \(Z\) component, the accuracy for the \(Y\) component is substantially better for the LP1 basis than for the \(\mathbf{L P} \mathbf{0}\) one. The difference in the accuracy results between the \(Y\) and \(Z\) components is likely explained by the fact that the function
\(x \mapsto z_{i}(x)\) is rather flat, so it is much better approximated by LP0 basis functions than \(x \mapsto y_{i}(x)\). The GPU speed-up reaches 65 .
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#C & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \text { av }}\) & \(M S E_{Z, \text { av }}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 3 & 125 & -4.021483 & -4.131725 & -0.900286 & 0.24 & 0.18 & 2.28 \\
\hline 0.1 & 5 & 500 & -4.290881 & -4.695769 & -1.551480 & 22.04 & 15.15 & 3.46 \\
\hline 0.05 & 7 & 2000 & -4.541253 & -5.022405 & -2.281332 & 1161.68 & 789.92 & 23.22 \\
\hline 0.02 & 10 & 12500 & -4.574551 & -5.143310 & -3.228237 & 170855.89 & 113131.26 & 2614.86 \\
\hline
\end{tabular}

LP1 local polynomials, \(d=4\), \#C \(=[3 \sqrt{d \sqrt{N}}-5], M=(d+1) N^{2}\).

Next, results for \(d=6\) are shown. Thus, we compare Table 8 below with Table 5 . Again, for a given precision on the \(Z\) component of the solution, we observe substantial improvements in the sequential and vectorized codes, but no such gains on the GPU version. In contrast, the accuracy of the \(Y\) approximation is, as in the \(d=4\) case, substantially better. Moreover, whereas we were not able to do computations for \(\Delta_{t}=0.02\) with the LP0 basis due to insufficient GPU memory, we are now able to make these calculations with the LP1 basis. The GPU speed-up reaches 75 , which is lower than the LP0 basis speed-up factor, as expected.
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#С & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 175 & -3.504153 & -3.668801 & -0.461077 & 0.39 & 0.25 & 2.45 \\
\hline 0.1 & 3 & 700 & -3.804091 & -3.911488 & -1.133263 & 56.61 & 37.39 & 5.10 \\
\hline 0.05 & 4 & 2800 & -4.075928 & -4.231639 & -1.791519 & 4422.97 & 2795.09 & 58.60 \\
\hline 0.02 & 6 & 17500 & -3.809734 & -4.529827 & -2.689432 & \(*\) & \(*\) & 43745.92 \\
\hline
\end{tabular}

In the high dimensional \(d=11\) setting shown in Table 9, we compare with Table 6. We observe a speed-up of order 125 compared to the sequential implementation.
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#C & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \text { av }}\) & \(M S E_{Z, \text { av }}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 2000 & -3.271648 & -3.368051 & -1.455388 & 290.84 & 178.24 & 8.05 \\
\hline 0.2 & 3 & 4000 & -3.269004 & -3.403994 & -1.975300 & 51043.88 & 32251.24 & 408.04 \\
\hline
\end{tabular}

TABLE 9
LP1 local polynomials, \(d=11\).

In the remainder of this section, we present results in dimension \(d=12\) to \(d=16\) (in Tables \(10,11,12,13\) and 14 , respectively) for which the capacity of the GPU is maximally used to provide the highest possible accuracy. The GPU speed-up reaches up to 118 compared to the sequential implementation.

\section*{Appendix.}
A.1. Proof of Proposition 2.1. It is known from [9, Proposition 3.1] that it is sufficient to show that there is a continuous \(C_{\rho}: \mathbb{R} \rightarrow[1, \infty)\) such that, for all \(\Lambda \geq 0\), \(\lambda \in[0, \Lambda]\), and \(y \in \mathbb{R}^{d}\),
\[
\begin{equation*}
\frac{p_{\operatorname{logis.}( }^{(\mu)}(y)}{C_{\rho}(\Lambda)} \leq \int_{\mathbb{R}^{d}} p_{\text {logis. }}^{(\mu)}(y+z \sqrt{\lambda}) \exp \left(-\frac{|z|^{2}}{2}\right) \mathrm{d} z \leq C_{\rho}(\Lambda) p_{\text {logis. }}^{(\mu)}(y) \tag{A.1}
\end{equation*}
\]
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \(\# \mathrm{C}\) & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \text { av }}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 2000 & -3.111153 & -3.232051 & -1.297737 & 646.55 & 394.63 & 10.03 \\
\hline 0.2 & 3 & 4000 & -3.214096 & -3.272644 & -1.821935 & 174473.26 & 108795.96 & 2086.94 \\
\hline
\end{tabular}

LP1 local polynomials, \(d=12\).
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \(\# \mathrm{C}\) & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 3000 & -2.995413 & -3.153302 & -1.460911 & 2129.23 & 1304.44 & 18.55 \\
\hline 0.2 & 2 & 4000 & -3.022855 & -3.158471 & -1.649632 & 2854.72 & 1749.02 & 24.24 \\
\hline
\end{tabular}

LP1 local polynomials, \(d=13\).

The proof is given for \(d=1\); generalization to higher dimensions is obvious because the multidimensional density is just the product of the one-dimensional densities over the components. Moreover, for simplicity the proof is given for \(\mu=1\), as generality in this parameter does not change the proof. For simplicity, we will write \(p_{\text {logis. }}^{(\mu)}(x)=p(x)\) in what follows.

In terms of the hyperbolic cosine function, the density can be expressed as
\[
p(x)=\frac{\exp (-x)}{(1+\exp (-x))^{2}}=\left(\exp \left(\frac{x}{2}\right)+\exp \left(-\frac{x}{2}\right)\right)^{-1}=\left(2 \cosh \left(\frac{x}{2}\right)\right)^{-1} .
\]

We define \(I(y, \lambda):=2 \int_{\mathbb{R}} p(y+z \sqrt{\lambda}) \exp \left(-\frac{z^{2}}{2}\right) \mathrm{d} z\), so that from the relation \(\cosh (x+\) \(y)=\cosh (x) \cosh (y)+\sinh (x) \sinh (y)\), we have that
\[
I(y, \lambda)=\int_{\mathbb{R}} \frac{\exp \left(-\frac{z^{2}}{2}\right)}{\cosh \left(\frac{y}{2}\right) \cosh \left(\frac{z \sqrt{\lambda}}{2}\right)+\sinh \left(\frac{y}{2}\right) \sinh \left(\frac{z \sqrt{\lambda}}{2}\right)} \mathrm{d} z:=I_{+}(y, \lambda)+I_{-}(y, \lambda)
\]
where \(I_{+,-}\)denotes respectively the integral on \(\mathbb{R}^{+}\)and \(\mathbb{R}^{-}\).
Upper bound. Suppose that \(y \geq 0\). Then, if \(z \geq 0\), it follows that \(\sinh (y / 2) \sinh (z \sqrt{\lambda} / 2) \geq\) 0 , whence
\[
I_{+}(y, \lambda) \leq \int_{\mathbb{R}_{+}} \frac{\exp \left(-\frac{z^{2}}{2}\right)}{\cosh \left(\frac{y}{2}\right) \cosh \left(\frac{z \sqrt{\lambda}}{2}\right)} \mathrm{d} z=2 \int_{\mathbb{R}_{+}} e^{-\frac{z^{2}}{2}} \mathrm{~d} z \times p(y)
\]

On the other hand, if \(z \leq 0\), then \(\sinh \left(\frac{y}{2}\right) \sinh \left(\frac{z \sqrt{\lambda}}{2}\right) \geq \cosh \left(\frac{y}{2}\right) \sinh \left(\frac{z \sqrt{\lambda}}{2}\right)\), therefore
\(I_{-}(y, \lambda) \leq \int_{\mathbb{R}_{-}} \frac{\exp \left(-\frac{z^{2}}{2}\right)}{\cosh \left(\frac{y}{2}\right)\left\{\cosh \left(\frac{z \sqrt{\lambda}}{2}\right)+\sinh \left(\frac{z \sqrt{\lambda}}{2}\right)\right\}} \mathrm{d} z=2 \int_{\mathbb{R}_{-}} \exp \left(-\frac{z^{2}}{2}-\frac{z \sqrt{\lambda}}{2}\right) \mathrm{d} z \times p(y)\).
Therefore, if \(y \geq 0\) then \(I(y, \lambda) \leq 2 \int_{\mathbb{R}} \exp \left(\frac{-z^{2}+(z)-\sqrt{\Lambda}}{2}\right) \mathrm{d} z \times p(y)\). Observing that \(I(y, \lambda)\) is symmetric in \(y\), thus the upper bound (A.1) is proved.

Lower bound. Suppose that \(y \geq 0\). For \(z \leq 0\), observe that \(\sinh \left(\frac{y}{2}\right) \sinh \left(\frac{z \sqrt{\lambda}}{2}\right) \leq 0\), whence
\[
I_{-}(y, \lambda) \geq \int_{\mathbb{R}_{-}} \frac{\exp \left(-\frac{z^{2}}{2}\right)}{\cosh \left(\frac{y}{2}\right) \cosh \left(\frac{z \sqrt{\lambda}}{2}\right)} \mathrm{d} z \geq 2 \int_{\mathbb{R}_{-}} \frac{\exp \left(-\frac{z^{2}}{2}\right)}{\cosh \left(\frac{z \sqrt{\Lambda}}{2}\right)} \mathrm{d} z \times p(y)
\]
\begin{tabular}{|c|c|c|c|c|c|c|c|c|}
\hline \(\Delta_{t}\) & \#C & M & \(M S E_{Y, \text { max }}\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline 0.2 & 2 & 2000 & -3.011673 & -3.092870 & -1.026128 & 3111.37 & 1915.67 & 27.36 \\
\hline 0.2 & 2 & 4000 & -3.029663 & \(-3.105833\) & -1.558935 & 6244.38 & 3904.64 & 80.26 \\
\hline
\end{tabular}

LP1 local polynomials, \(d=14\).
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \#С & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \mathrm{av}}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 5000 & -2.981181 & -3.106590 & -1.574532 & 17246.62 & 10522.20 & 226.37 \\
\hline
\end{tabular}

LP1 local polynomials, \(d=15\).

For \(z \geq 0\), we use that \(\sinh \left(\frac{y}{2}\right) \sinh \left(\frac{z \sqrt{\lambda}}{2}\right) \leq \cosh \left(\frac{y}{2}\right) \sinh \left(\frac{z \sqrt{\lambda}}{2}\right)\) to obtain
\[
\begin{aligned}
I_{+}(y, \lambda) & \geq \int_{\mathbb{R}_{+}} \frac{\exp \left(-\frac{z^{2}}{2}\right)}{\cosh \left(\frac{y}{2}\right)\left\{\cosh \left(\frac{z \sqrt{\lambda}}{2}\right)+\sinh \left(\frac{z \sqrt{\lambda}}{2}\right)\right\}} \mathrm{d} z \\
& \geq 2 \int_{\mathbb{R}_{+}} \exp \left(-\frac{z^{2}}{2}-\frac{z \sqrt{\Lambda}}{2}\right) \mathrm{d} z \times p(y)
\end{aligned}
\]

The result on \(y<0\) follows again from the symmetry of \(I(y, \lambda)\).
A.2. Stability results for discrete BSDE. We recall standard results borrowed to [11] and adapted to our setting, they are aimed at comparing two solutions of discrete BSDEs of the form (2.1) with different data. Namely, consider two discrete BSDEs, \(\left(Y_{1, i}, Z_{1, i}\right)_{0 \leq i<N}\) and \(\left(Y_{2, i}, Z_{2, i}\right)_{0 \leq i<N}\), given by
\[
\begin{aligned}
Y_{l, i} & =\mathbb{E}_{i}\left(g\left(X_{N}\right)+\sum_{j=i}^{N-1} f_{l, j}\left(X_{j}, Y_{l, j+1}, Z_{l, j}\right) \Delta_{t}\right) \\
\Delta_{t} Z_{l, i} & =\mathbb{E}_{i}\left(\left(g\left(X_{N}\right)+\sum_{j=i+1}^{N-1} f_{l, j}\left(X_{j}, Y_{l, j+1}, Z_{l, j}\right) \Delta_{t}\right) \Delta W_{i}\right),
\end{aligned}
\]
for \(i \in\{0, \ldots, N-1\}, l \in\{1,2\}\).
To allow the driver \(f_{1, i}\) to depend on the clouds of simulations (necessary in the analysis), we require that it is measurable w.r.t. \(\mathcal{F}_{T}\) instead of \(\mathcal{F}_{t_{i}}\) as usually.

Proposition A.1. Assume that \(\left(\mathbf{A}_{\mathbf{g}}\right)\) and \(\left(\mathbf{A}_{\mathbf{X}}\right)\) hold. Moreover, for each \(i \in\) \(\{0, \ldots, N-1\}\), assume that \(f_{1, i}\left(X_{i}, Y_{1, i+1}, Z_{1, i}\right) \in \mathbf{L}_{2}\left(\mathcal{F}_{T}\right)\) and \(f_{2}\) satisfies \(\left(\mathbf{A}_{\mathbf{f}}\right)\) with constants \(L_{f_{2}}\) and \(C_{f_{2}}\). Then, for any \(\gamma \in(0,+\infty)\) satisfying \(6 q\left(\Delta_{t}+\frac{1}{\gamma}\right) L_{f_{2}}^{2} \leq 1\), we have for \(0 \leq i<N\)
\[
\begin{aligned}
& \left|Y_{1, i}-Y_{2, i}\right|^{2} \Gamma_{i}+\sum_{j=i}^{N-1} \Delta_{t} \mathbb{E}_{i}\left(\left|Z_{1, j}-Z_{2, j}\right|^{2}\right) \Gamma_{j} \\
& \leq 3 C_{A .1}\left(\frac{1}{\gamma}+\Delta_{t}\right) \sum_{j=i}^{N-1} \Delta_{t} \mathbb{E}_{i}\left(\left|f_{1, j}\left(X_{j}, Y_{1, j+1}, Z_{1, j}\right)-f_{2, j}\left(X_{j}, Y_{1, j+1}, Z_{1, j}\right)\right|^{2}\right) \Gamma_{j}
\end{aligned}
\]
where \(\Gamma_{i}:=\left(1+\gamma \Delta_{t}\right)^{i}\) and \(C_{A .1}:=2 q+(1+T) e^{T / 2}\).
\begin{tabular}{|r|r|r||r|r|r|r|r|r|}
\hline\(\Delta_{t}\) & \(\# \mathrm{C}\) & \(M\) & \(M S E_{Y, \max }\) & \(M S E_{Y, \text { av }}\) & \(M S E_{Z, \mathrm{av}}\) & CPU & VE & GPU \\
\hline \hline 0.2 & 2 & 6000 & -2.795353 & -2.959375 & -1.588716 & 45587.17 & 27507.47 & 669.28 \\
\hline
\end{tabular}

LP1 local polynomials, \(d=16\).

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