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# Adaptive force biasing algorithms: new convergence results and tensor approximations of the bias 

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#### Abstract

A modification of the Adaptive Biasing Force method is introduced, in which the free energy is approximated by a sum of tensor products of one-dimensional functions. This enables to handle a larger number of reaction coordinates than the classical algorithm. We prove the algorithm is well-defined and prove the long-time convergence toward a regularized version of the free energy for an idealized version of the algorithm. Numerical experiments demonstrate that the method is able to capture correlations between reaction coordinates.


## 1 Introduction

Consider $x \in \mathbb{T}^{D}$ a vector representing the positions of particles with periodic boundary conditions $(\mathbb{T}=\mathbb{R} / \mathbb{Z})$, and a potential energy $V \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$. We are interested in computing expectations of the form

$$
\frac{1}{\int_{\mathbb{T}^{D}} e^{-\beta V(x)} \mathrm{d} x} \int_{\mathbb{T}^{D}} \varphi(x) e^{-\beta V(x)} \mathrm{d} x=: \int_{\mathbb{T}^{D}} \varphi \mathrm{~d} \mu_{V, \beta}
$$

where $\varphi: \mathbb{T}^{D} \rightarrow \mathbb{R}$ is called an observable and $\mathrm{d} \mu_{V, \beta}:=e^{-\beta V(x)} \mathrm{d} x$ is the Gibbs law with potential $V$ and inverse temperature $\beta>0$. The large dimension $D$ is so significant that, in practice, these quantities have to be computed with Markov Chain Monte Carlo (MCMC) algorithms, which consist in approximating the average of $\varphi$ with respect to $\mu_{V, \beta}$ along dynamics that are ergodic with respect to $\mu_{V, \beta}$. A typical sampler is the overdamped Langevin dynamics

$$
\mathrm{d} X_{t}=-\nabla V\left(X_{t}\right) \mathrm{d} t+\sqrt{2 \beta^{-1}} \mathrm{~d} B_{t}
$$

where $\left(B_{t}\right)_{t \geqslant 0}$ is a Brownian motion over $\mathbb{T}^{D}$. It is ergodic with invariant measure $\mu_{V, \beta}$, so that

$$
\frac{1}{t} \int_{0}^{t} \varphi\left(X_{s}\right) \mathrm{d} s \underset{t \rightarrow \infty}{\longrightarrow} \int_{\mathbb{T}^{D}} \varphi \mathrm{~d} \mu_{V, \beta} .
$$

Nevertheless, the convergence of the process (or, in practice, of any alternative Markov process with invariant measure $\mu_{V, \beta}$ ) toward its equilibrium in the long-time limit may be very
slow. This is due to the so-called metastability phenomenon, according to which the process remains for long times in some region of the space, with very rare transitions from one of these metastable regions to another. This is related to the multi-modality of the Gibbs measure and the fact MCMC algorithms typically perform local moves, so that leaving a mode of the target measure $\mu_{V, \beta}$ is a rare event. We refer to [18] for more details on this topic. For this reason, several adaptive methods have been developed in order to force the process to leave the metastable traps faster. Among those, we focus on the adaptive biasing force (ABF) algorithm, which may be seen as a particular Importance Sampling method. The general idea is to run a biased process

$$
\begin{equation*}
\mathrm{d} \tilde{X}_{t}=-\nabla V\left(\tilde{X}_{t}\right) \mathrm{d} t+\nabla V_{\text {bias }, t}\left(\tilde{X}_{t}\right)+\sqrt{2 \beta^{-1}} \mathrm{~d} B_{t} \tag{1}
\end{equation*}
$$

where the biasing potential $V_{\text {bias }, t}$ is adaptively constructed from the past trajectory $\left(\tilde{X}_{s}\right)_{s \in[0, t]}$ in such a way that it is expected to converge to some $V_{\text {bias }, \infty}$. Expectations with respect to $\mu_{V, \beta}$ are then recovered through a reweighting step, assuming that ergodicity still holds:

Classically, in such an Importance Sampling scheme, the aim is to design a target bias $V_{\text {bias, } \infty}$ such that two conditions are met: 1) sampling the biased equilibrium $\mu_{V-V_{\text {bias }, \infty}, \beta}$ is simpler than the initial problem (i.e. the corresponding overdamped Langevin process is less metastable) and 2) the biased equilibrium is not too far from the initial target so that the exponential weights in (2) do not cause the asymptotical variance of the estimator to skyrocket.

In the ABF algorithm, this issue is addressed with the use of so-called reaction coordinates (or collective variables) and the associated free energy as a bias. Reaction coordinates consist of a small number $d \ll D$ of macroscopic coordinates of the whole microscopic system $x \in \mathbb{T}^{D}$. These coordinates are defined through a $\operatorname{map} \xi: \mathbb{T}^{D} \rightarrow \mathcal{M}$ where $\mathcal{M}$ is a manifold of dimension $d$. In molecular dynamics, for example, $x \in \mathbb{T}^{D}$ is a vector which gathers the positions of all the different atoms of the system of interest, and $\xi(x)$ typically represents some distances between particular pairs of atoms, or angles formed by some triplets of atoms. These reaction coordinates should be chosen to capture the main causes of the metastability of the system. More precisely, $\xi\left(X_{t}\right)$ should converge to equilibrium as slowly as $X_{t}$, while the conditional laws $\mathcal{L}(X \mid \xi(X)=z)$ for fixed $z \in \mathcal{M}$ when $X \sim \mu_{V, \beta}$ should be easier to sample (see [19] for more detailed considerations). In other words, $\xi(x)$ should be a low-dimensional representation of $x$ that captures the slow variables of the system.

To these reaction coordinates $\xi$ is associated the corresponding free energy $A: \mathcal{M} \rightarrow \mathbb{R}$, given by

$$
A(z)=-\frac{1}{\beta} \ln \int_{\left\{x \in \mathbb{T}^{D}, \xi(x)=z\right\}} e^{-\beta V(x)} \delta_{\xi(x)-z}(\mathrm{~d} x)
$$

where $\delta_{\xi(x)-z}$ is the so-called delta measure, which can be defined from the Lebesgue measure on the submanifold $\left\{x \in \mathbb{T}^{D}, \xi(x)=z\right\}$ through the co-area formula, see for example [20, Section 3.2.1]. This definition ensures that, if $X$ is a random variable with law $\mu_{V, \beta}$ on $\mathbb{T}^{D}$, then $\xi(X)$ is a random variable with law $\mu_{A, \beta}$ on $\mathcal{M}$. The heuristic of the ABF algorithm is the following. Suppose that $\mathcal{M}$ is compact. If we were to sample from the process

$$
\mathrm{d} Y_{t}=-\nabla(V-A \circ \xi)\left(Y_{t}\right) \mathrm{d} t+\sqrt{2 \beta^{-1}} \mathrm{~d} B_{t}
$$

the equilibrium would be $\mu_{V-A \circ \xi, \beta}$, whose image through $\xi$, by definition of $A$, is the uniform measure on $\mathcal{M}$. This means that there would be no more metastability along $\xi$, since all
the regions of $\mathcal{M}$ would be equally visited by $\xi\left(Y_{t}\right) \mathrm{s}$. Unfortunately, it is not possible to use directly this free-energy biased dynamics in practice, since it would require the knowledge of $A$ and thus the computation of expectations in large dimension. The idea of the ABF method is to learn $A$ on the fly, i.e. to run a process $\left(\tilde{X}_{t}\right)_{t \geqslant 0}$ solving (1) with a biasing potential $V_{\text {bias }, t}$ constructed from $\left(\tilde{X}_{s}\right)_{s \in[0, t]}$ and designed to target $A \circ \xi$ in the longtime limit.

In practice, the choice of good reaction coordinates is a difficult problem. Up to recently, their definition has been based on the knowledge and intuition of specialists. The question of the automatic learning of suitable reaction coordinates is currently a vivid research area, see for instance [7, 8, Moreover, some techniques like the orthogonal space random walk [22] provide a general way to construct new reaction coordinates from previous ones. Due to these recent progresses, one would like to consider a relatively large $d$. In ABF, $V_{\text {bias, }}$ is a function of the $d$ reaction coordinates. From a numerical point of view, since $V_{b i a s, t}$ is adaptively learned on the fly, its values have to be kept in memory, which requires a grid whose size typically scales exponentially with $d$. This limits the application of ABF to small dimensional reaction coordinates $(d \leqslant 4)$. The aim of the present work is to lift this limitation by approximating $V_{\text {bias,t }}$ using a sum of tensor products of one-dimensional functions, which reduces the size of the memory to $\mathcal{O}(d m)$ where $m$ is the number of tensor terms. Remark that this can in turn help for the definition of good reaction coordinates, by considering as candidates a relatively large number of reaction coordinates and then conduct a statistical study to select or combine some of them. Nevertheless, this question exceeds the scope of the present work, in which $\xi$ is supposed to be given.

Note that the question of increasing the number of reaction coordinates in adaptive biasing algorithms has also been considered in the Bias-Exchange algorithm introduced in [26], where several replicas of the system are run in parallel, each associated with a one-dimensional reaction coordinate. The replicas exchange their bias according to some Metropolis-Hastings probability, so that each replica eventually feels the bias in all the different directions of the reaction coordinates. Nevertheless, in this case where one-dimensional reaction coordinates are treated independently one from the others, the system remains very sensitive to correlations between reaction coordinates (the same goes for the generalized ABF introduced in [27]), contrary to the algorithm introduced in the present work.

Besides, let us mention that numerical methods involving both tensor approximation and Monte Carlo methods for molecular dynamics are also introduced in [16, 24] for other purposes.

In the rest of this introduction we provide a presentation of the ABF algorithm we consider in this work in a simple framework, and refer to Section 4 for generalizations. The presentation is divided into two parts. In Section 1.1, we present the reference ABF algorithm we consider, without the tensor-product approximation. In Section 1.2, we introduce the tensor-product approximation of the bias. These two ingredients are then combined to yield the TensorABF algorithm in Section 1.3. The two algorithms and associated convergence proofs of the reference ABF algorithm and of the tensor-product approximation are presented separately since we think they have their own interest.

### 1.1 Free energy and the ABF algorithm

Let us first present the ABF algorithm in a simple framework (see [15, 10, 19 , for more general settings). From now on, we write

$$
\mu=\mu_{V, \beta}
$$

seen both as a probability law and as the density of the latter with respect to the Lebesgue measure.

Let us assume that $\mathcal{M}=\mathbb{T}^{d}$ and that, for all $x=(q, z) \in \mathbb{T}^{D}=\mathbb{T}^{p} \times \mathbb{T}^{d}, \xi(x)=\xi(q, z)=z$ where $p=D-d$.

At first sight, this may seem a very restrictive choice of reaction coordinates. But, using extended variables (see [13]), this can be applied actually in very general contexts. We refer the reader to Section 4 for more details on this point.

The associated free energy for $z \in \mathbb{T}^{d}$ is then

$$
A(z)=-\frac{1}{\beta} \ln \int_{\mathbb{T}^{p}} e^{-\beta V(q, z)} \mathrm{d} q
$$

Following the previous discussion, our aim is then to define for all time $t \geqslant 0$ a function $A_{t}$ on $\mathbb{T}^{d}$ and to sample the process

$$
\left\{\begin{align*}
\mathrm{d} Q_{t} & =-\nabla_{q} V\left(Q_{t}, Z_{t}\right) \mathrm{d} t+\sqrt{2 \beta^{-1}} \mathrm{~d} B_{t}^{1}  \tag{3}\\
\mathrm{~d} Z_{t} & =-\nabla_{z} V\left(Q_{t}, Z_{t}\right) \mathrm{d} t+\nabla_{z} A_{t}\left(Z_{t}\right) \mathrm{d} t+\sqrt{2 \beta^{-1}} \mathrm{~d} B_{t}^{2}
\end{align*}\right.
$$

where $B^{1}$ and $B^{2}$ are independent Brownian motions respectively of dimension $p$ and $d$, in such a way that $A_{t}$ gets close to $A$ in large time.

Note that the free energy $A$ satisfies

$$
\nabla_{z} A(z)=\frac{\int_{\mathbb{T}^{p}} \nabla_{z} V(q, z) e^{-\beta V(q, z)} \mathrm{d} q}{\int_{\mathbb{T}^{p}} e^{-\beta V(q, z)} \mathrm{d} q}=\mathbb{E}_{\mu}\left[\nabla_{z} V(Q, Z) \mid Z=z\right]
$$

The following alternative equivalent characterization of $A$ will be useful in the sequel. Let us define

$$
\begin{equation*}
H:=\left\{f \in H^{1}\left(\mathbb{T}^{d}\right): \int_{\mathbb{T}^{d}} f(z) \mathrm{d} z=0\right\} \tag{4}
\end{equation*}
$$

and let us denote by $\mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$ the set of probability measures on $\mathbb{T}^{p} \times \mathbb{T}^{d}=\mathbb{T}^{D}$. For all $\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$ and $f \in H$, let us define

$$
\mathcal{E}_{\nu}(f):=\int_{\mathbb{T}^{p} \times \mathbb{T}^{d}}\left|\nabla_{z} V(q, z)-\nabla_{z} f(z)\right|^{2} \mathrm{~d} \nu(q, z) .
$$

It holds that, up to an additive constant (like the potential $V$, the free energy is in fact always defined up to an additive constant), $A$ is the unique minimizer in $H$ of the functional $\mathcal{E}_{\mu}$, i.e.

$$
\begin{equation*}
A=\underset{f \in H}{\operatorname{argmin}} \mathcal{E}_{\mu}(f) . \tag{5}
\end{equation*}
$$

At time $t \geqslant 0$, a trajectory $\left(Q_{s}, Z_{s}\right)_{s \in[0, t]}$ of (3) is available. Let $\nu_{t}$ be the probability measure on $\mathbb{T}^{p} \times \mathbb{T}^{d}=\mathbb{T}^{D}$ defined as follows: for all $\varphi \in \mathcal{C}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$,

$$
\begin{equation*}
\int_{\mathbb{T}^{p} \times \mathbb{T}^{d}} \varphi \mathrm{~d} \nu_{t}=\left(\int_{0}^{t} e^{-\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s\right)^{-1} \int_{0}^{t} \varphi\left(Q_{s}, Z_{s}\right) e^{-\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s \tag{6}
\end{equation*}
$$

We call $\nu_{t}$ the unbiased occupation distribution of the process. By the ergodic limit (2), $\nu_{t}$ is expected to converge to $\mu$ as $t$ goes to infinity (at least if $A_{t}$ does not change too fast with $t$ ).

However, note that $\nu_{t}$ is a singular probability measure, so that the minimization problem

$$
\inf _{f \in H} \mathcal{E}_{\nu_{t}}(f)
$$

is ill-posed. To circumvent this difficulty, one may consider two different alternatives to regularize the problem which we detail hereafter. Consider a smooth symmetric positive density kernel $K \in \mathcal{C}^{\infty}\left(\mathbb{T}^{d} \times \mathbb{T}^{d}, \mathbb{R}_{+}\right)$with

$$
\begin{equation*}
\int_{\mathbb{T}^{d}} K(y, z) \mathrm{d} z=1 \quad \text { and } \quad K(y, z)=K(z, y) \quad \forall y, z \in \mathbb{T}^{d} \tag{7}
\end{equation*}
$$

In practice, $K(y, \cdot)$ should be close to a Dirac mass at $y$ (see Theorem 2 below). For instance, a possible choice for $K$ would be the so-called von-Mises kernel for a given small parameter $\varepsilon>0$, i.e.

$$
\begin{equation*}
K(y, z) \propto \prod_{i=1}^{d} \exp \left(-\frac{1}{\varepsilon^{2} / 2} \sin ^{2}\left(\frac{z_{i}-y_{i}}{2}\right)\right) \tag{8}
\end{equation*}
$$

Now, consider also a regularization parameter $\lambda \geqslant 0$. For all $\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$ and all $f \in H$, we define

$$
\begin{equation*}
\mathcal{J}_{\nu}(f):=\int_{\mathbb{T}^{p} \times \mathbb{T}^{d} \times \mathbb{T}^{d}}\left|\nabla_{y} V(q, y)-\nabla_{z} f(z)\right|^{2} K(y, z) \mathrm{d} z \mathrm{~d} \nu(q, y)+\lambda \int_{\mathbb{T}^{d}}\left|\nabla_{z} f(z)\right|^{2} \mathrm{~d} z \tag{9}
\end{equation*}
$$

Note that, as $K(y, \cdot)$ converges toward the Dirac mass at $y$ and $\lambda$ goes to 0 , for all $f \in H, \mathcal{J}_{\nu}(f)$ converges towards $\mathcal{E}_{\nu}(f)$. The interest of introducing $\mathcal{J}_{\nu}$ is that, thanks to the regularization, the minimization problem is now well-posed:

Proposition 1. Assume that either $K>0$ on $\mathbb{T}^{d} \times \mathbb{T}^{d}$ or $\lambda>0$. Then, for all $\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$, $\mathcal{J}_{\nu}$ admits a unique minimizer in $H$.

In summary, in the whole article, we work under the following conditions.
Assumption 1. $V \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right), D \geqslant 3, \beta>0, \lambda \geqslant 0$ and $K \in \mathcal{C}^{\infty}\left(\mathbb{T}^{d} \times \mathbb{T}^{d}, \mathbb{R}_{+}\right)$satisfies (7). Moreover, either $K>0$ or $\lambda>0$.

We now have all the elements to define the reference ABF algorithm in this work, see Algorithm 1 below.

Remark that, contrary to the cases studied in other theoretical works like [19, 1, 4], in Algorithm 1, the bias $A_{t}$ is piecewise constant in time, with updates at the times $t_{k}, k \in$ $\llbracket 1, N_{u p} \rrbracket$. This is due to the fact that, as will be detailed in Section 1.2, the bias updates are numerically demanding in our case, and thus we cannot perform them at each timestep.

We prove in Section 2 the long-time convergence of Algorithm 1;
Theorem 2. Under Assumption 1, let $\left(Q_{t}, Z_{t}, A_{t}\right)_{t \geqslant 0}$ be given by Algorithm 1 (with $N_{u p}=$ $+\infty)$. Then, as $t \rightarrow+\infty$, almost surely, $\nu_{t}$ weakly converges toward $\mu$ and

$$
\left\|\nabla A_{t}-\nabla A_{*}\right\|_{\infty} \underset{t \rightarrow \infty}{\longrightarrow} 0
$$

where $A_{*}$ is the unique minimizer in $H$ of $\mathcal{J}_{\mu}$. Moreover, $A_{*}$ satisfies

$$
\begin{align*}
\int_{\mathbb{T}^{d}}\left|\nabla A(z)-\nabla A_{*}(z)\right|^{2} & \left(\int_{\mathbb{T}^{p} \times \mathbb{T}^{d}} K(y, z) \mu(q, y) d q d y\right) d z \\
& \leqslant 4\left\|\nabla^{2} A\right\|_{\infty} \sup _{y \in \mathbb{T}^{d}} \int_{\mathbb{T}^{d}}|y-z|^{2} K(y, z) d z+2 \lambda \int_{\mathbb{T}^{d}}|\nabla A(z)|^{2} d z \tag{10}
\end{align*}
$$

```
Algorithm 1 ABF algorithm
    Input:
    Initial condition \(\left(q_{0}, z_{0}\right) \in \mathbb{T}^{p} \times \mathbb{T}^{d}\)
    Brownian motion \(\left(B_{t}^{1}, B_{t}^{2}\right)_{t \geqslant 0}\) on \(\mathbb{T}^{p} \times \mathbb{T}^{d}\)
    Regularization parameters \(K, \lambda\)
    Update period \(T_{u p}>0\), number of updates \(N_{u p} \in \mathbb{N}_{*}\), total simulation time \(T_{t o t}=T_{u p} N_{u p}\)
    Output:
    Estimated free energy \(A_{T_{\text {tot }}} \in H\)
    Trajectory \(\left(Q_{t}, Z_{t}\right)_{t \in\left[0, T_{\text {tot }}\right]} \in \mathcal{C}\left(\left[0, T_{\text {tot }}\right], \mathbb{T}^{p} \times \mathbb{T}^{d}\right)\)
    Begin:
    Set \(\left(Q_{0}, Z_{0}\right)=\left(q_{0}, z_{0}\right)\).
    Set \(A_{0}(z)=0\) for all \(z \in \mathbb{T}^{d}\).
    Set \(t_{k}=k T_{u p}\) for all \(k \in \llbracket 0, N_{u p} \rrbracket\).
    for \(k \in \llbracket 1, N_{u p} \rrbracket\) do
        Set \(A_{t}=A_{t_{k-1}}\) for all \(t \in\left[t_{k-1}, t_{k}\right)\).
        Set \(\left(Q_{t}, Z_{t}\right)_{t \in\left[t_{k-1}, t_{k}\right]}\) to be the solution of (3) with initial condition \(\left(Q_{t_{k-1}}, Z_{t_{k-1}}\right)\) at
    time \(t_{k-1}\).
        Set \(A_{t_{k}}\) to be the minimizer in \(H\) of \(\mathcal{J}_{\nu_{t_{k}}}\) given by (6) and (9).
    end for
    Return \(A_{T_{\text {tot }}}\) and \(\left(Q_{t}, Z_{t}\right)_{t \in\left[0, T_{\text {tot }}\right]}\).
```

Note that (10) implies that, as $\lambda$ and $\sup _{y \in \mathbb{T}^{d}} \int_{\mathbb{T}^{d}}|y-z|^{2} K(y, z) \mathrm{d} z$ go to zero, $A_{*}$ converges in $H$ to $A$ (which corresponds to $\lambda=0$ and $K(y, z)=\delta_{y}(z)$ ).

The long-time convergence of a similar ABF algorithm has been established in [1] but in the case where the empirical measure $\nu_{t}$ is replaced by the law of the process at time $t$. Rather than a self-interacting process (i.e. a single trajectory with memory), this corresponds to a system of $N$ interacting particles (with no memory), and more precisely to the mean-field limit as $N$ goes to infinity of this system. Moreover, a result similar to Theorem 2 has been established in [4] for a closely related self-interacting process, the adaptive biasing potential algorithm. In addition, in the recent preprint [5], a similar result is established for the ABF algorithm but when the occupation measure is not unbiased (see the discussion in Section 4.3).

### 1.2 Tensor approximation

This section focuses on the minimization step of Algorithm 1. Assumption 1 is enforced. Fix $\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$. For all $f \in H$, the cost function $\mathcal{J}_{\nu}(f)$ defined by (9) is equal to

$$
\begin{equation*}
\mathcal{J}_{\nu}(f)=C_{\nu}+(1+\lambda) \int_{\mathbb{T}^{d}}\left|F_{\nu}(z)-\nabla f(z)\right|^{2} \theta_{\nu}(z) \mathrm{d} z \tag{11}
\end{equation*}
$$

with some constant $C_{\nu}$ independent from $f$ and where, for all $z \in \mathbb{T}^{d}$,

$$
\begin{align*}
\theta_{\nu}(z) & :=\frac{1}{\lambda+1}\left(\lambda+\int_{\mathbb{T}^{p} \times \mathbb{T}^{d}} K(y, z) \mathrm{d} \nu(q, y)\right)  \tag{12}\\
F_{\nu}(z) & :=\frac{1}{(\lambda+1) \theta_{\nu}(z)} \int_{\mathbb{T}^{p} \times \mathbb{T}^{d}} \nabla_{y} V(q, y) K(y, z) \mathrm{d} \nu(q, y) . \tag{13}
\end{align*}
$$

Note that, under Assumption 1, $\theta_{\nu}$ is the density of a probability measure, bounded from below by $(\lambda+\min K) /(1+\lambda)>0$. Moreover, since $K$ is smooth and bounded, so are $\theta_{\nu}$ and $F_{\nu}$. Note that neither the additive constant $C_{\nu}$ nor the multiplication by $1+\lambda$ affect the problem of minimizing $\mathcal{J}_{\nu}$. As a consequence, the unique minimizer $f_{*}$ of $\mathcal{J}_{\nu}$ on $H$ (see Proposition (1) is equivalently the unique minimizer of

$$
\begin{equation*}
H \ni f \mapsto \widetilde{\mathcal{J}}_{\nu}(f):=\int_{\mathbb{T}^{d}}\left|F_{\nu}(z)-\nabla f(z)\right|^{2} \theta_{\nu}(z) \mathrm{d} z \tag{14}
\end{equation*}
$$

The gradient of $\widetilde{\mathcal{J}}_{\nu}$ at $f_{*}$ is the Helmholtz projection in $L^{2}\left(\theta_{\nu}\right)$ of $F_{\nu}$. The Euler-Lagrange equation associated to the minimization problem of $\widetilde{\mathcal{J}}_{\nu}$ over $H$ is

$$
\begin{equation*}
\nabla \cdot\left(\theta_{\nu}\left(\nabla f_{*}-F_{\nu}\right)\right)=0 \tag{15}
\end{equation*}
$$

where $\nabla$. denotes the divergence operator. When $d$ is small ( $d=2$ in [1]), as $t$ increases, the functions $\theta_{\nu_{t}}$ and $\theta_{\nu_{t}} F_{\nu_{t}}$ are updated and kept in memory on a discrete grid of dimension $M^{d}$ for some $M \in \mathbb{N}^{*}$, and the Euler equation is solved with standard PDE techniques. However this is not sustainable if one wants to consider a larger number of reaction coordinates. For this reason, we now present a method to approximate $f_{*}$ by a sum of tensor products, namely by a function $f_{m} \in H$ which reads as follows

$$
\forall z:=\left(z_{1}, \cdots, z_{d}\right) \in \mathbb{T}^{d}, \quad f_{m}(z)=\sum_{k=1}^{m} \prod_{j=1}^{d} r_{k, j}\left(z_{j}\right)
$$

for some $m \in \mathbb{N}_{*}$ and some functions $r_{k, j}: \mathbb{T} \rightarrow \mathbb{R}$ for $1 \leq j \leq d$ and $1 \leq k \leq m$. See [14] for a general overview on tensor methods.

Let $g$ be a simple tensor product function, i.e. a function such that for all $z=\left(z_{1}, \cdots, z_{d}\right) \in$ $\mathbb{T}^{d}, g(z)=\prod_{j=1}^{d} r_{j}\left(z_{j}\right)$ for some $r_{1}, \cdots, r_{d} \in H^{1}(\mathbb{T})$. Such a simple tensor product function will be denoted herefater by $g=\bigotimes_{j=1}^{d} r_{j}$.

If $g$ belongs to $H$, its (Lebesgue) integral vanishes, which is equivalent to the fact there exists $i \in \llbracket 1, d \rrbracket$ such that the (Lebesgue) integral of $r_{i}$ vanishes. This motives the introduction of the following subspaces of $H$ : for $i \in \llbracket 1, d \rrbracket$, define

$$
\Sigma_{i}:=\left\{g \in H, g=\bigotimes_{j=1}^{d} r_{j} \text { with } r_{j} \in H^{1}(\mathbb{T}) \text { for all } j \in \llbracket 1, d \rrbracket \text { and } \int_{\mathbb{T}} r_{i}\left(z_{i}\right) \mathrm{d} z_{i}=0\right\} .
$$

Proposition 3. Under Assumption 1, for all $\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right), i \in \llbracket 1, d \rrbracket$ and $f \in H$, there always exists at least one minimizer in $\Sigma_{i}$ to the optimization problem

$$
\begin{equation*}
\min _{g \in \Sigma_{i}} \mathcal{J}_{\nu}(f+g) \tag{16}
\end{equation*}
$$

From Proposition 3, the greedy algorithm described in Algorithm 2 below is well-defined.
In Section 3 is established the following:
Theorem 4. Under Assumption 1, let $f_{*}$ be the minimizer of $\mathcal{J}_{\nu}$ in $H$ and $f_{m}=\operatorname{Greedy}\left(\nu, f_{0}, m\right)$ as given by Algorithm (2) for some $\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$, $f_{0} \in H$ and $m \in \mathbb{N}$. Then

$$
\left\|f_{m}-f_{*}\right\|_{H^{1}} \underset{m \rightarrow+\infty}{\longrightarrow} 0
$$

```
Algorithm \(2 \operatorname{Greedy}\left(\nu, f_{0}, m\right)\)
    Input:
    Probability measure \(\nu \in \mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)\)
    Initial guess \(f_{0} \in H\)
    number of tensor terms \(m \in \mathbb{N}_{*}\)
    Output:
    \(f_{m} \in H\).
    Begin:
    \(n=0\)
    while \(n<m\) do
        for \(i \in \llbracket 1, d \rrbracket\) do
            Find \(g_{n}:=\bigotimes_{j=1}^{d} r_{n, j}\) a minimizer of \(g \mapsto \mathcal{J}_{\nu}\left(f_{n}+g\right)\) over \(g \in \Sigma_{i}\) (i.e. with
    \(r_{n, j} \in H^{1}(\mathbb{T})\) for all \(1 \leq j \leq d\) and \(\left.\int_{\mathbb{T}} r_{n, i}=0\right)\)
            Set \(f_{n+1}=f_{n}+g_{n}\).
            Increment \(n \leftarrow n+1\).
        end for
    end while
    Return \(f_{m}\).
```

The interest of Algorithm 2 is that at each iteration, one only has to compute $d$ onedimensional functions, which makes it possible to implement even if $d$ is relatively large (say $4<d<10$ ). Notice that the price to pay when going from the original problem of minimizing $\mathcal{J}_{\nu}$ over $H$ to the problem (16) is that the Euler-Lagrange equations associated to the initial problem are linear (since $\mathcal{J}_{\nu}$ is a quadratic functional) whereas the Euler-Lagrange equations associated to (16) are nonlinear. This is due to the fact that the quadratic functional is minimized over a non-linear space in (16).

In practice, the problem of minimizing $\mathcal{J}_{\nu}(f+g)$ over $g=\bigotimes_{j=1}^{d} r_{j} \in \Sigma_{i}$ is solved through the Alternating Least Square method [12], which is a fixed point procedure on the EulerLagrange equation (15): the $r_{j}$ 's are optimized one after the other, the others being fixed, repeatedly. This amounts to solving a system of one-dimensional elliptic PDEs of the form

$$
\begin{equation*}
\partial_{z_{j}}\left(a_{j} \partial_{z_{j}} r_{j}\right)\left(z_{j}\right)-b_{j}\left(z_{j}\right) r_{j}\left(z_{j}\right)=c_{j}\left(z_{j}\right) \tag{17}
\end{equation*}
$$

with

$$
\begin{aligned}
& a_{j}\left(z_{j}\right)=\int_{\mathbb{T}^{d-1}}\left(\prod_{l \neq j} r_{l}\left(z_{l}\right)\right)^{2} \theta_{\nu}(z) \mathrm{d} z_{\neq j} \\
& b_{j}\left(z_{j}\right)=\sum_{h \neq j} \int_{\mathbb{T}^{d-1}}\left|\partial_{z_{h}} \prod_{l \neq j} r_{l}\left(z_{l}\right)\right|^{2} \theta_{\nu}(z) \mathrm{d} z_{\neq j} \\
& c_{j}\left(z_{j}\right)=\int_{\mathbb{T}^{d-1}}\left(\prod_{l \neq j} r_{l}\left(z_{l}\right)\right) \partial_{z_{j}}\left(F_{\nu, j} \theta_{\nu}\right)(z) \mathrm{d} z_{\neq j}-\sum_{h \neq j} \int_{\mathbb{T}^{d-1}} \partial_{z_{h}}\left(\prod_{l \neq j} r_{l}\left(z_{l}\right)\right) F_{\nu, h}(z) \theta(z) \mathrm{d} z_{\neq j},
\end{aligned}
$$

where $\mathrm{d} z_{\neq j}$ means that all variables except the $j^{\text {th }}$ are integrated and $F_{\nu, j}$ denotes the $j^{\text {th }}$ component of $F_{\nu}$. If for all $y=\left(y_{1}, \cdots, y_{d}\right), z=\left(z_{1}, \cdots, z_{d}\right) \in \mathbb{T}^{d}, K(y, z)=\prod_{i=1}^{d} K_{i}\left(y_{i}, z_{i}\right)$ for
some functions $K_{i}: \mathbb{T} \times \mathbb{T} \rightarrow \mathbb{R}$ for all $1 \leq i \leq d$ (like the kernel (8)), for $\nu=\nu_{t}$ given by (6),

$$
(1+\lambda) a_{j}\left(z_{j}\right)=\prod_{l \neq j}\left\|r_{l}\right\|_{L^{2}(\mathbb{T})}^{2}+\int_{\mathbb{T}^{D}}\left(\prod_{l \neq j} \int_{\mathbb{T}} r_{l}^{2}\left(z_{l}\right) K_{l}\left(y_{l}, z_{l}\right) \mathrm{d} z_{l}\right) K_{j}\left(y_{j}, z_{j}\right) \mathrm{d} \nu_{t}(q, y),
$$

which can be computed without computing $F_{\nu_{t}}(z)$ and $\theta_{\nu_{t}}(z)$ for all $z \in \mathbb{T}^{d}$ (which would be improssible in practice). The same holds for $b_{j}$ and $c_{j}$.

### 1.3 The tensor ABF algorithm

As already explained above, the main objective of this work is to introduce a new algorithm to adapt the standard ABF approach to multi-dimensional reaction coordinates. Combining Algorithms 1 and 2, the Tensor ABF (TABF) algorithm is described in Algorithm 3 below. Note that, for the sake of clarity, it has been kept relatively simple. In particular, we haven't addressed here the question of time and space discretization.

Moreover, the proofs of convergence of Algorithm 1 and Algorithm 2 also have their own interest. The convergence of Algorithm 1 is based on the so-called ordinary differential equation method [6], and requires specific contractivity bounds. The convergence of Algorithm 2 is an adaptation of the proof of convergence of greedy algorithms [9, the main difficulty being to deal with the zero average constraint in $H$.

The paper is organized as follows. Section 2 is devoted to the proof of Theorem 2. In Section 3, we prove Theorem 4. Section 4 provides a detailed discussion on practical considerations and possible variations of the algorithm. Finally, some numerical experiments with the TABF algorithm are reported in Section 5.

```
Algorithm 3 TABF algorithm
    Input:
    Initial condition \(\left(q_{0}, z_{0}\right) \in \mathbb{T}^{p} \times \mathbb{T}^{d}\)
    Brownian motion \(\left(B_{t}^{1}, B_{t}^{2}\right)_{t \geqslant 0}\) on \(\mathbb{T}^{p} \times \mathbb{T}^{d}\)
    Regularization parameters \(K, \lambda\)
    Update period \(T_{u p}>0\), number of updates \(N_{u p} \in \mathbb{N}_{*}\), total simulation time \(T_{t o t}=T_{u p} N_{u p}\)
    Number of tensor terms by update \(m \in \mathbb{N}_{*}\)
    Output:
    Estimated free energy \(A_{T_{\text {tot }}} \in H\)
    Trajectory \(\left(Q_{t}, Z_{t}\right)_{t \in\left[0, T_{\text {tot }}\right]} \in \mathcal{C}\left(\left[0, T_{\text {tot }}\right], \mathbb{T}^{p} \times \mathbb{T}^{d}\right)\)
    Begin:
    Set \(\left(Q_{0}, Z_{0}\right)=\left(q_{0}, z_{0}\right)\).
    Set \(A_{0}(z)=0\) for all \(z \in \mathbb{T}^{d}\).
    Set \(t_{k}=k T_{u p}\) for all \(k \in \llbracket 0, N_{u p} \rrbracket\).
    for \(k \in \llbracket 1, N_{u p} \rrbracket\) do
        Set \(A_{t}=A_{t_{k-1}}\) for all \(t \in\left[t_{k-1}, t_{k}\right)\).
        Set \(\left(Q_{t}, Z_{t}\right)_{t \in\left[t_{k-1}, t_{k}\right]}\) to be the solution of (3) with value \(\left(Q_{t_{k-1}}, Z_{t_{k-1}}\right)\) at time \(t_{k-1}\).
        Set \(f_{m}=\operatorname{Greedy}\left(\nu_{t_{k}}, A_{t_{k-1}}, m\right)\) given by Algorithm 2 where \(\nu_{t_{k}}\) is given by (6).
        Set \(A_{t_{k}}=f_{m}\).
    Return \(A_{T_{t o t}}\) and \(\left(Q_{t}, Z_{t}\right)_{t \in\left[0, T_{t o t}\right]}\).
```


## 2 Proof of the long-time convergence

In the whole Section 2 we consider the ABF process $\left(Q_{t}, Z_{t}, A_{t}\right)_{t \geqslant 0}$ obtaind through Algorithm 1 (with $N_{u p}=+\infty$ ), and Assumption 1 holds.
Lemma 5. For all $r \in \mathbb{N}^{*}$ and all multi-index $\alpha \in \mathbb{N}^{d}$, there exists a constant $C_{\alpha}>0$ such that, for all $t \geqslant 0,\left\|\partial^{\alpha} A_{t}\right\|_{\infty} \leqslant C_{\alpha}$.

Proof. Since $\mathbb{R}_{+} \ni t \mapsto A_{t}$ is piecewise constant, we may assume that $t=t_{k}=k T_{u p}$ for some $k \in \mathbb{N}$ without loss of generality. Using the notation of Section 1.2, $A_{t}$ is then the minimizer over $H$ of $\widetilde{\mathcal{J}}_{\nu_{t}}$ defined in (14). Recall that for all $f \in H$,

$$
\widetilde{\mathcal{J}}_{\nu_{t}}(f)=\int_{\mathbb{T}^{d}}\left|F_{\nu_{t}}(z)-\nabla f(z)\right|^{2} \theta_{\nu_{t}}(z) \mathrm{d} z
$$

Remark that $\theta_{\nu_{t}}$ is bounded from below uniformly in $t$ and $z$ by $(\lambda+\min K) /(1+\lambda)>0$, and similarly all the derivatives in $z$ of $\theta_{\nu_{t}}$ and of $F_{\nu_{t}}$ are bounded in $L^{\infty}\left(\mathbb{T}^{d}\right)$ by constants which depend on $K$ and $V$ but not on $t$. The Euler-Lagrange equation associated to the minimization of $\widetilde{\mathcal{J}}_{\nu_{t}}$ reads

$$
\begin{equation*}
\nabla \cdot\left(\theta_{\nu_{t}} \nabla A_{t}\right)=\nabla \cdot\left(\theta_{\nu_{t}} F_{\nu_{t}}\right) \tag{18}
\end{equation*}
$$

By elliptic regularity (cf. [2]), $A_{t}$ is thus $\mathcal{C}^{\infty}$ and, differentiating (18), multiplying it by derivatives of $\nabla A_{t}$ and integrating, we classically get by induction that

$$
\int_{\mathbb{T}^{d}}\left|\partial^{\alpha} \nabla A_{t}\right|^{2} \theta_{\nu_{t}} \leqslant C_{\alpha}
$$

where $\alpha \in \mathbb{N}^{d}$ is any multi-index, for some constant $C_{\alpha}>0$ which does not depend on $t$. Conclusion follows from Sobolev embeddings.

Theorem 2 will be a direct corollary of:
Proposition 6. Almost surely, $\nu_{t} \underset{t \rightarrow \infty}{\text { weak }} \mu$.
The proof of Proposition 6 is postponed to the end of this section. Let us prove that indeed, given the latter, Theorem 2 holds:

Proof of Theorem 2. By the arguments of the previous proof, for all $t \geq 0$, the function $\mathbb{T}^{d} \ni z \mapsto \theta_{\nu_{t}}(z)$ is bounded and Lipschitz with constants which are uniform in $t$. Hence, for any $\varepsilon>0$, we can find $N_{\varepsilon} \in \mathbb{N}^{*}$ and a finite set of points $z_{1}, \cdots, z_{N_{\varepsilon}} \in \mathbb{T}^{d}$ such that for all $z \in \mathbb{T}^{d}$, there exists $i_{z} \in \llbracket 1, N_{\varepsilon} \rrbracket$ such that, for all $t>0,\left|\theta_{\nu_{t}}\left(z_{i_{z}}\right)-\theta_{\nu_{t}}(z)\right| \leqslant \varepsilon$. The same holds for $\theta_{\mu}$. On the other hand, according to Proposition 6, almost surely,

$$
\sup _{i \in \llbracket 1, N_{\varepsilon} \rrbracket}\left|\theta_{\nu_{t}}\left(z_{i}\right)-\theta_{\mu}\left(z_{i}\right)\right| \underset{t \rightarrow \infty}{\longrightarrow} 0,
$$

so that $\left\|\theta_{\nu_{t}}-\theta_{\mu}\right\|_{\infty}$ goes to zero as $t \rightarrow \infty$. Similar arguments enable us to obtain the same results for all the derivatives of $\theta_{\nu_{t}}$ and for $F_{\nu_{t}}$ and all its derivatives. Note that $A_{*}$ is the minimizer of

$$
\widetilde{\mathcal{J}}_{\mu}(f)=\int_{\mathbb{T}^{d}}\left|F_{\mu}(z)-\nabla f(z)\right|^{2} \theta_{\mu}(z) \mathrm{d} z
$$

Let $t=t_{n}$ for some $n \in \mathbb{N}^{*}$. The associated Euler-Lagrange equations associated with the two minimization problems on $A_{t}$ and $A_{*}$ lead to

$$
\begin{equation*}
\nabla \cdot\left(\theta_{\mu} \nabla\left(A_{t}-A_{*}\right)\right)=\nabla \cdot\left(\theta_{\nu_{t}} F_{\nu_{t}}-\theta_{\mu} F_{\mu}-\left(\theta_{\nu_{t}}-\theta_{\mu}\right) \nabla A_{t}\right) \tag{19}
\end{equation*}
$$

Multiplying this equality by $A_{t}-A_{*}$, integrating and using the uniform control on $\nabla A_{t}$ established in Lemma 5 (and the lower bound on $\theta_{\mu}$ ), we get that

$$
\int_{\mathbb{T}^{d}}\left|\nabla\left(A_{t}(z)-A_{*}(z)\right)\right|^{2} \mathrm{~d} z \underset{t \rightarrow \infty}{\longrightarrow} 0 .
$$

More generally, differentiating (19), multiplying it by derivatives of $A_{t}-A_{*}$, integrating and using the uniform controls of the derivatives of $A_{t}$, we obtain by induction that

$$
\int_{\mathbb{T}^{d}}\left|\nabla \partial^{\alpha}\left(A_{t}-A_{*}\right)(z)\right|^{2} \mathrm{~d} z \underset{t \rightarrow \infty}{\longrightarrow} 0
$$

for all multi-index $\alpha \in \mathbb{N}^{d}$. The first statement of Theorem 2 then follows from Sobolev embeddings.

Finally, inequality (10) stems from the fact that $\mathcal{J}_{\mu}\left(A_{*}\right) \leqslant \mathcal{J}_{\mu}(A)$. More precisely, using that

$$
\int_{\mathbb{T}^{d}} \nabla A(y) \mu(q, y) \mathrm{d} q=\int_{\mathbb{T}^{d}} \nabla_{y} V(q, y) \mu(q, y) \mathrm{d} q
$$

we get that for all $f \in H, \mathcal{J}_{\mu}(f)=\widehat{\mathcal{J}}_{\mu}(f)+\int\left|\nabla_{y} V\right|^{2} \mathrm{~d} \mu-\int|\nabla A|^{2} \mathrm{~d} \mu$ where

$$
\widehat{\mathcal{J}}_{\mu}(f)=\int_{\mathbb{T}^{p} \times \mathbb{T}^{d} \times \mathbb{T}^{d}}|\nabla A(y)-\nabla f(z)|^{2} K(y, z) \mathrm{d} z \mathrm{~d} \mu(q, y)+\lambda \int_{\mathbb{T}^{d}}|\nabla f(z)|^{2} \mathrm{~d} z
$$

In other words, $\mathcal{J}_{\mu}$ and $\widehat{\mathcal{J}}_{\mu}$ only differ by an additive constant, so that $A_{*}$ is the minimizer of $\widehat{\mathcal{J}}_{\mu}$ over $H$. Then

$$
\begin{aligned}
& \int_{\mathbb{T}^{p} \times \mathbb{T}^{d} \times \mathbb{T}^{d}}\left|\nabla A(z)-\nabla A_{*}(z)\right|^{2} K(y, z) \mathrm{d} z \mathrm{~d} \mu(q, y) \\
& \leqslant 2 \widehat{\mathcal{J}}_{\mu}\left(A_{*}\right)+2 \int_{\mathbb{T}^{p} \times \mathbb{T}^{d} \times \mathbb{T}^{d}}|\nabla A(y)-\nabla A(z)|^{2} K(y, z) \mathrm{d} z \mathrm{~d} \mu(q, y) \\
& \leqslant 2 \widehat{\mathcal{J}}_{\mu}(A)+2 \int_{\mathbb{T}^{p} \times \mathbb{T}^{d} \times \mathbb{T}^{d}}|\nabla A(y)-\nabla A(z)|^{2} K(y, z) \mathrm{d} z \mathrm{~d} \mu(q, y) \\
& \leqslant 2 \lambda \int_{\mathbb{T}^{d}}|\nabla A(z)|^{2} \mathrm{~d} z+4 \int_{\mathbb{T}^{p} \times \mathbb{T}^{d} \times \mathbb{T}^{d}}|\nabla A(y)-\nabla A(z)|^{2} K(y, z) \mathrm{d} z \mathrm{~d} \mu(q, y) \\
& \leqslant 2 \lambda \int_{\mathbb{T}^{d}}|\nabla A(z)|^{2} \mathrm{~d} z+4\left\|\nabla^{2} A\right\|_{\infty}^{2} \sup _{y \in \mathbb{T}^{d}} \int_{\mathbb{T}^{d}}|y-z|^{2} K(y, z) \mathrm{d} z
\end{aligned}
$$

The rest of the section is dedicated to the proof of Proposition 6, which follows the so-called ordinary differential equation (ODE) method [6].

### 2.1 Time change and the ODE method

Following an idea of [4], we introduce the (random) time change:

$$
\tau(t):=\int_{0}^{t} e^{-\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s
$$

so that

$$
\nu_{t}=\frac{1}{\tau(t)} \int_{0}^{t} \delta_{Q_{s}, Z_{s}} \tau^{\prime}(s) \mathrm{d} s=\frac{1}{\tau(t)} \int_{0}^{\tau(t)} \delta_{Q_{\tau^{-1}(s)}, Z_{\tau^{-1}(s)}} \mathrm{d} s .
$$

In other words, considering the time-changed process $\bar{X}_{t}:=\left(Q_{\tau^{-1}(t)}, Z_{\tau^{-1}(t)}\right)$ and its occupation measure

$$
\begin{equation*}
\bar{\nu}_{t}=\frac{1}{t} \int_{0}^{t} \delta_{\bar{X}_{s}} \mathrm{~d} s, \tag{20}
\end{equation*}
$$

then $\nu_{t}=\bar{\nu}_{\tau(t)}$. Since, at a fixed time $t \geqslant 0, A_{t}$ is smooth and with Lebesgue integral zero, there always exists $z \in \mathbb{T}^{d}$ such that $A_{t}(z)=0$, so that

$$
\begin{equation*}
\left\|A_{t}\right\|_{\infty} \leqslant \sqrt{d} / 2\left\|\nabla A_{t}\right\|_{\infty} \tag{21}
\end{equation*}
$$

where we used that $\sqrt{d} / 2$ is the diameter of $\mathbb{T}^{d}$. Together with Lemma 5 , this implies that in particular, $\tau(t)$ goes to infinity with $t$.

Denoting $S_{t}(x):=A_{\tau^{-1}(t)}(z)$ for all $x=(q, z) \in \mathbb{T}^{p} \times \mathbb{T}^{d}$, the inhomogeneous Markov process $\bar{X}$ solves the SDE

$$
\begin{equation*}
\mathrm{d} \bar{X}_{t}=-e^{\beta S_{t}\left(\bar{X}_{t}\right)} \nabla\left(V-S_{t}\right)\left(\bar{X}_{t}\right)+\sqrt{2 \beta^{-1} e^{\beta S_{t}\left(\bar{X}_{t}\right)}} \mathrm{d} \bar{B}_{t} \tag{22}
\end{equation*}
$$

where $\left(\bar{B}_{t}\right)_{t \geqslant 0}$ is a standard Brownian motion on $\mathbb{T}^{D}$, obtained from $\left(B_{t}\right)_{t \geqslant 0}$ through rescaling. We denote by $\left(L_{t}\right)_{t \geqslant 0}$ its infinitesimal generator, defined by: for all $\varphi \in \mathcal{C}^{2}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)=\mathcal{C}^{2}\left(\mathbb{T}^{D}\right)$ and all $x \in \mathbb{T}^{p} \times \mathbb{T}^{d}$,

$$
L_{t} \varphi(x)=\lim _{h \rightarrow 0} \frac{\mathbb{E}\left(\varphi\left(\bar{X}_{t+h}\right) \mid \bar{X}_{t}=x\right)-\varphi(x)}{h}
$$

whenever the limit exists. Here,

$$
L_{t} \varphi(x)=\left(-\nabla\left(V-S_{t}\right)(x) \cdot \nabla \varphi(x)+\frac{1}{\beta} \Delta \varphi(x)\right) e^{\beta S_{t}(x)}
$$

We denote by $\left(P_{s}^{(t)}\right)_{s \geqslant 0}$ the Markov semi-group generated by $L_{t}$ for a fixed $t$. Formally, $P_{s}^{(t)}=e^{s L_{t}}$. For all $t \geqslant 0$ the unique invariant measure of $L_{t}$ is $\mu$ (see [4, Proposition 3.1]). From Lemma 5 and the bound (21), we consider $C_{0}>0$ such that $S_{t} \in \mathcal{B}_{C_{0}}$ for all $t \geqslant 0$ where

$$
\mathcal{B}_{C_{0}}:=\left\{S \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right), \int_{\mathbb{T}^{D}} S(x) \mathrm{d} x=0,\|S\|_{\mathcal{C}^{2}\left(\mathbb{T}^{D}\right)} \leqslant C_{0}\right\}
$$

The principle of the ODE method is the following: for large values of the time $t$, the evolution of $\bar{\nu}_{t}$ is slow (because of the $t^{-1}$ factor in (20)). Hence, for $1 \ll s \ll t$, in principle, it holds that $S_{u} \simeq S_{t}$ for $u \in[t, t+s]$, so that

$$
\begin{equation*}
\bar{\nu}_{t+s}=\frac{t}{t+s} \bar{\nu}_{t}+\frac{s}{s+t}\left(\frac{1}{s} \int_{t}^{t+s} \delta_{\bar{X}_{u}} \mathrm{~d} u\right) \simeq \frac{t}{t+s} \bar{\nu}_{t}+\frac{s}{s+t} \mu . \tag{23}
\end{equation*}
$$

In other words, the evolution of $\bar{\nu}_{t}$ approximately follows the deterministic flow

$$
\partial_{t} m_{t}=\frac{1}{t}\left(\mu-m_{t}\right)
$$

which converges to $\mu$, so that $\bar{\nu}_{t}$ (hence $\nu_{t}$ ) should also converge to $\mu$.
Following [6] (to which we refer for more details), we will make rigorous this heuristic. For $t, s>0$, set

$$
\begin{equation*}
\varepsilon_{t}(s)=\int_{e^{t}}^{e^{t+s}} \frac{\delta_{\bar{X}_{u}}-\mu}{u} \mathrm{~d} u \tag{24}
\end{equation*}
$$

It is said that $\mathbb{R}_{+} \ni t \mapsto \bar{\nu}_{e^{t}}$ is an asymptotic pseudotrajectory of the flow

$$
\begin{equation*}
\partial_{t} m_{t}=\left(\mu-m_{t}\right) \tag{25}
\end{equation*}
$$

if, for all $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ and $T>0$,

$$
\lim _{t \rightarrow \infty} \sup _{s \in[0, T]}\left|\varepsilon_{t}(s) \varphi\right|=0
$$

(cf [6, Proposition 3.5]). Note that (25) admits $\mu$ as a unique attractor, in the sense that

$$
d_{T V}\left(m_{t}, \mu\right) \leqslant e^{-t} d_{T V}\left(m_{0}, \mu\right) \underset{t \rightarrow \infty}{\longrightarrow} 0
$$

for all $m_{0} \in \mathcal{P}\left(\mathbb{T}^{D}\right)$, where

$$
\begin{equation*}
d_{T V}\left(m_{1}, m_{2}\right)=\sup _{A \in \mathcal{B}\left(\mathbb{T}^{D}\right)}\left|m_{1}(A)-m_{2}(A)\right| \tag{26}
\end{equation*}
$$

is the total variation distance on $\mathcal{P}\left(\mathbb{T}^{D}\right)$, whose topology is stronger than the weak convergence. In $(26)$, we have denoted by $\mathcal{B}\left(\mathbb{T}^{D}\right)$ the set of borelian subsets of $\mathbb{T}^{D}$.

As a consequence, according to [6, Theorem 3.7], the limit set (i.e. the set of (weak) limits of extracted sequences) of an asymptotic pseudotrajectory of the flow (25) is necessarily reduced to the singleton $\{\mu\}$. In other words, by compactness (since $\mathbb{T}^{D}$ is compact, so is $\mathcal{P}\left(\mathbb{T}^{D}\right)$ ), an asymptotic pseudotrajectory of (25) necessarily converges to $\mu$ (in the weak sense).

As a conclusion, Proposition 6 is a corollary of:
Proposition 7. Almost surely, $\mathbb{R}_{+} \ni t \mapsto \bar{\nu}_{e^{t}}$ is an asymptotic pseudotrajectory of (25).
In order to prove this result, we will need some quantitative estimates. Indeed, note that, for the approximation (23) to hold, the speed of convergence of $P^{(t)}$ toward $\mu$ should be uniform in $t$. As we will see below, this is a direct consequence of the estimate of Lemma 5 . The approximation also requires that the time evolution of $A_{t}$ is controlled in some sense.

### 2.2 Preliminary estimates

For a fixed $S \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$, consider $L_{S}$ defined for $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ by

$$
L_{S} \varphi(x)=\left(-\nabla(V-S)(x) \cdot \nabla \varphi(x)+\frac{1}{\beta} \Delta \varphi(x)\right) e^{\beta S(x)}
$$

which is the infinitesimal generator of the SDE

$$
\mathrm{d} X_{t}^{S}=-e^{\beta S\left(X_{t}^{S}\right)} \nabla(V-S)\left(X_{t}^{S}\right) \mathrm{d} t+\sqrt{2 \beta^{-1} e^{\beta S\left(X_{t}^{S}\right)}} \mathrm{d} B_{t}
$$

Denote $\left(P_{t}^{S}\right)_{t \geqslant 0}$ the associated (homogeneous) semi-group and $\Gamma_{S}$ the associated carré-duchamps operator, defined for $\varphi, \psi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ and all $x \in \mathbb{T}^{D}$ by

$$
\Gamma_{S}(\varphi, \psi)(x):=\frac{1}{2}\left(L_{S}(\varphi \psi)-\varphi L_{S} \psi-\psi L_{S} \varphi\right)(x)=\beta^{-1} e^{\beta S(x)} \nabla \varphi(x) \cdot \nabla \psi(x)
$$

and $\Gamma_{S}(\varphi):=\Gamma_{S}(\varphi, \varphi)$. By classical elliptic regularity arguments, if $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ then $P_{t}^{S} \varphi \in$ $\mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$, in particular $\mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ is a core for $L_{S}$, see [3, Section 1.13]. More precisely each derivative of $P_{t}^{S} \varphi$ is uniformly bounded over all finite time interval, which ensures the validity
of the computations in the proofs of the next lemmas. Integrating twice by parts, it can be easily seen that for all $\varphi, \psi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\int_{\mathbb{T}^{D}} \varphi(x) L_{S} \psi(x) \mu(\mathrm{d} x)=\int_{\mathbb{T}^{D}} \psi(x) L_{S} \varphi(x) \mu(\mathrm{d} x)
$$

in other words $L_{S}$ is a self-adjoint operator on $L^{2}(\mu)$.
Lemma 8. Let us assume that $D \geq 3$. Then, there exists $C_{1}>0$ such that for all $S \in \mathcal{B}_{C_{0}}$, $\left(\mu, L_{S}\right)$ satisfies a Poincaré inequality and a Sobolev inequality both with constant $C_{1}$, in the sense that for all $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\begin{aligned}
\|\varphi\|_{L^{2}(\mu)}^{2} & \leqslant C_{1} \int_{\mathbb{T}^{D}} \Gamma_{S}(\varphi) d \mu \\
\|\varphi\|_{L^{p}(\mu)}^{2} & \leqslant C_{1}\left(\|\varphi\|_{L^{2}(\mu)}^{2}+\int_{\mathbb{T}^{D}} \Gamma_{S}(\varphi) d \mu\right),
\end{aligned}
$$

where $p=\frac{2 D}{D-2}$.
Proof. For $S=0$, the first inequality is the classical Poincaré inequality, which holds here since the density of $\mu$ with respect to the Lebesgue measure is bounded above and below away from zero, see [3, Proposition 5.1.6]. As a consequence, there exists $c>0$ such that for all $S \in \mathcal{B}_{C_{0}}$ and $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\|\varphi\|_{L^{2}(\mu)}^{2} \leqslant c \int_{\mathbb{T}^{D}}|\nabla \varphi|^{2} \mathrm{~d} \mu \leqslant c e^{\beta C_{0}} \int_{\mathbb{T}^{D}} \Gamma_{S}(\varphi) \mathrm{d} \mu .
$$

Similarly, from the Sobolev inequality satisfied by the Lebesgue measure on $\mathbb{T}^{D}$ [3, Section 6],

$$
\begin{aligned}
\|\varphi\|_{L^{p}(\mu)}^{2} & \leqslant\|\mu\|_{\infty}^{2 / p}\|\varphi\|_{L^{p}\left(\mathbb{T}^{D}\right)}^{2} \\
& \leqslant C\|\mu\|_{\infty}^{2 / p}\left(\|\varphi\|_{L^{2}\left(\mathbb{T}^{D}\right)}^{2}+\|\nabla \varphi\|_{L^{2}\left(\mathbb{T}^{D}\right)}^{2}\right) \\
& \leqslant C\|\mu\|_{\infty}^{2 / p}\left\|\mu^{-1}\right\|_{\infty}^{2}\left(\|\varphi\|_{L^{2}(\mu)}^{2}+\|\nabla \varphi\|_{L^{2}(\mu)}^{2}\right) \\
& \leqslant C e^{\beta C_{0}}\|\mu\|_{\infty}^{2 / p}\left\|\mu^{-1}\right\|_{\infty}^{2}\left(\|\varphi\|_{L^{2}(\mu)}^{2}+\int_{\mathbb{T}^{D}} \Gamma_{S}(\varphi) \mathrm{d} \mu\right) .
\end{aligned}
$$

These inequalities, in turn, yield the following estimates:
Lemma 9. There exist $C_{2}>0$ such that, for all $S \in \mathcal{B}_{C_{0}}, t \geqslant 0$ and $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\begin{aligned}
\left\|P_{t}^{S} \Pi \varphi\right\|_{L^{2}(\mu)} & \leqslant e^{-t / C_{2}}\|\Pi \varphi\|_{L^{2}(\mu)} \\
\left\|P_{t}^{S} \varphi\right\|_{\infty} & \leqslant \frac{C_{2}}{\min \left(1, t^{d / 2}\right)}\|\varphi\|_{L^{2}(\mu)} \\
\left\|\nabla P_{t}^{S} \varphi\right\|_{\infty} & \leqslant \frac{C_{2}}{\min (1, \sqrt{t})}\|\varphi\|_{\infty}
\end{aligned}
$$

with $\Pi \varphi:=\varphi-\int_{\mathbb{T}^{D}} \varphi d \mu$.

Proof. The first estimate is a usual consequence of the Poincaré inequality, see [3, Proposition 5.1.3]. The second one, namely the ultracontractivity of the semi-group, is a consequence of the Sobolev inequality (see [3, Theorem 6.3.1]). The last one can be established thanks to the Bakry-Emery calculus (see [3, Section 1.16] for an introduction), by showing that $L_{S}$ satisfies a curvature estimate, as we now detail. We would like to compare $\left|\nabla P_{t} \varphi\right|^{2}$ and $P_{t}\left(\varphi^{2}\right)$. A seminal idea of the Bakry-Emery calculus is that quantities of the form $\Theta\left(P_{t} \varphi\right)$ and $P_{t} \Theta(\varphi)$, where $\Theta$ is some operator can be linked through the interpolation $P_{t-s} \Theta\left(P_{s} \varphi\right), s \in[0, t]$, so that $\Theta\left(P_{t} \varphi\right)-P_{t} \Theta(\varphi)=\int_{0}^{t} \partial_{s}\left(P_{t-s} \Theta\left(P_{s} \varphi\right)\right) \mathrm{d} s$. When differentiating with respect to $s$, we obtain quantities of the form $-2 P_{t-s} \Gamma_{\Theta}\left(P_{s} \varphi\right)$ for some operator $\Gamma_{\Theta}$, which is of a form similar to the interpolation ( $\Theta$ being replaced by $\Gamma_{\Theta}$ ).

More precisely, when $\Theta(\varphi)=\varphi^{2}$, then $\Gamma_{\Theta}$ is the usual carré-du-champ operator, and when $\Theta(\varphi)=|\nabla \varphi|^{2}$ we end up with

$$
\Gamma_{\nabla, S}(\varphi)=\frac{1}{2} L_{S}\left(|\nabla \varphi|^{2}\right)-\nabla \varphi \cdot \nabla L_{S} \varphi,
$$

for $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$. Writing $[\varphi, \psi]=\varphi \psi-\psi \varphi$, we compute

$$
\begin{aligned}
\Gamma_{\nabla, S}(\varphi)= & \sum_{i=1}^{D}\left(\Gamma_{S}\left(\partial_{x_{i}} \varphi\right)+\partial_{x_{i}} \varphi\left[\partial_{x_{i}}, L_{S}\right] \varphi\right) \\
\geqslant & \sum_{i=1}^{D}\left[\beta^{-1} e^{-\beta\|S\|_{\infty}}\left|\nabla \partial_{x_{i}} \varphi\right|^{2}-\beta^{-1} e^{\beta\|S\|_{\infty}}\left|\nabla \partial_{x_{i}}(V-S) \| \nabla \varphi\right|\left|\partial_{x_{i}} \varphi\right|\right. \\
& \left.-\beta\left|\partial_{x_{i}} S\right| e^{\beta\|S\|_{\infty}}\left|\partial_{x_{i}} \varphi \| \nabla(V-S) \cdot \nabla \varphi+\frac{1}{\beta} \Delta \varphi\right|\right] \\
\geqslant & -c|\nabla \varphi|^{2}
\end{aligned}
$$

for some $c>0$ which is uniform over $S \in \mathcal{B}_{C_{0}}$. Now, following [23, Lemma 4], we want to consider the interpolation between $\alpha(t)\left|\nabla P_{t} \varphi\right|^{2}+\left(P_{t} \varphi\right)^{2}$ and $\alpha(0) P_{t}|\nabla \varphi|^{2}+P_{t}\left(\varphi^{2}\right)$ for some $\alpha$ with $\alpha(0)=0<\alpha(t)$. For fixed $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right), x \in \mathbb{T}^{D}$ and $t \geqslant 0$, we set for all $s \in[0, t]$

$$
\Psi(s)=\alpha(s) P_{t-s}^{S}\left|\nabla P_{s}^{S} \varphi\right|^{2}(x)+e^{\beta C_{0}} P_{t-s}^{S}\left(P_{s}^{S} \varphi\right)^{2}(x)
$$

with $\alpha(s)=(1-\exp (-2 c t)) / c$, so that

$$
\begin{aligned}
\partial_{s} \Psi(s) & =P_{t-s}^{S}\left(-2 \alpha(s) \Gamma_{\nabla, S}+\alpha^{\prime}(s)|\nabla \cdot|^{2}-2 e^{\beta C_{0}} \Gamma_{S}\right)\left(P_{s}^{S} \varphi\right)(x) \\
& \leqslant\left(2 \alpha(s) c+\alpha^{\prime}(s)-2\right) P_{t-s}^{S}\left|\nabla P_{s}^{S} \varphi\right|^{2}(x)=0 .
\end{aligned}
$$

In particular,

$$
\alpha(t)\left|\nabla P_{t}^{S} \varphi\right|^{2}(x) \leqslant \Psi(t) \leqslant \Psi(0)=e^{\beta C_{0}} P_{t}^{S} \varphi^{2}(x) \leqslant e^{\beta C_{0}}\|\varphi\|_{\infty}^{2}
$$

which yields the desired estimate.
Lemma 10. There exists $C_{3}>0$ such that for all $S \in \mathcal{B}_{C_{0}}$, the operator $R_{S}$ defined for all $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ by

$$
R_{S} \varphi=-\int_{0}^{\infty} P_{t}^{S} \Pi \varphi d t
$$

satisfies $L_{S} R_{S}=R_{S} L_{S}=\Pi$ and, for all $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\begin{equation*}
\left\|R_{S} \varphi\right\|_{\infty}+\left\|\nabla R_{S} \varphi\right\|_{\infty}+\left\|\Delta R_{S} \varphi\right\|_{\infty} \leqslant C_{3}\|\varphi\|_{\infty} \tag{27}
\end{equation*}
$$

Proof. We follow the proof of [6, Section 5.2 and Lemma 5.1]. First, from Lemma 9 (and using the fact that $\left\|P_{t}^{S} \varphi\right\|_{\infty} \leqslant\|\varphi\|_{\infty}$ for all $t \geqslant 0$ ),

$$
\begin{aligned}
\int_{0}^{\infty}\left\|P_{t}^{S} \Pi \varphi\right\|_{\infty} d t & \leqslant \int_{0}^{1}\|\Pi \varphi\|_{\infty} d t+\int_{1}^{\infty}\left\|P_{t}^{S} \Pi \varphi\right\|_{\infty} d t \\
& \leqslant 2\|\varphi\|_{\infty}+C_{2} \int_{1}^{\infty}\left\|P_{t-1}^{S} \Pi \varphi\right\|_{L^{2}(\mu)} d t \\
& \leqslant 2\|\varphi\|_{\infty}+C_{2} \int_{1}^{\infty} e^{-(t-1) / C_{2}}\|\Pi \varphi\|_{L^{2}(\mu)} d t \\
& \leqslant 2\left(1+C_{2}^{2}\right)\|\varphi\|_{\infty}
\end{aligned}
$$

and similarly, using the fact that $\left\|\nabla P_{t}^{S} \Pi \varphi\right\|_{\infty}=\left\|\nabla P_{1} P_{t-1}^{S} \Pi \varphi\right\|_{\infty} \leq C_{2}\left\|P_{t-1}^{S} \Pi \phi\right\|_{\infty}$ for $t \geqslant 1$,

$$
\begin{aligned}
\int_{0}^{\infty}\left\|\nabla P_{t}^{S} \Pi \varphi\right\|_{\infty} d t & \leqslant \int_{0}^{1} \frac{C_{2}}{\sqrt{t}}\|\Pi \varphi\|_{\infty} d t+C_{2} \int_{1}^{\infty}\left\|P_{t-1}^{S} \Pi \varphi\right\|_{\infty} d t \\
& \leqslant C_{2}\left(2+2\left(1+C_{2}^{2}\right)\right)\|\varphi\|_{\infty}
\end{aligned}
$$

In particular $R_{S} \varphi$ and $\nabla R_{S} \varphi$ are well defined in $L^{\infty}\left(\mathbb{T}^{D}\right)$ for $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$. Moreover, using the fact that, from Lemma 9. $\left\|P_{t}^{S} \Pi \varphi\right\|_{\infty} \leqslant C_{2} e^{-(t-1) / C_{2}}\|\Pi \varphi\|_{L^{2}(\mu)} \rightarrow 0$ as $t \rightarrow+\infty$,

$$
\begin{aligned}
L_{S} R_{S} \varphi & =-\int_{0}^{\infty} L_{S} P_{t}^{S} \Pi \varphi d t \\
& =-\int_{0}^{\infty} \partial_{t}\left(P_{t}^{S} \Pi \varphi\right) d t=\Pi \varphi
\end{aligned}
$$

The case of $R_{S} L_{S}$ is similar: since $\mu$ is invariant for $L_{S}, L_{S} \Pi=L_{S}=\Pi L_{S}$, and thus $P_{t}^{S} \Pi L_{S}=$ $P_{t}^{S} L_{S} \Pi=\partial_{t}\left(P_{t}^{S} \Pi \varphi\right)$ for all $t \geqslant 0$.

As a consequence,

$$
\left|\Delta R_{S} \varphi\right| \leqslant e^{\beta C_{0}}\left|e^{\beta S} \Delta R_{S} \varphi\right| \leqslant e^{\beta C_{0}}\left(\left\|e^{\beta S} \nabla(V-S) \cdot \nabla R_{S} \varphi\right\|_{\infty}+\|\Pi \varphi\|_{\infty}\right) \leqslant C_{3}\|\varphi\|_{\infty}
$$

for some $C_{3}>0$ uniform over $S \in \mathcal{B}_{C_{0}}$, which yields the desired result.
Lemma 11. There exist $C_{4}>0$ such that for all $S_{1}, S_{2} \in \mathcal{B}_{C_{0}}$ and $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\left\|R_{S_{1}} \varphi-R_{S_{2}} \varphi\right\|_{\infty} \leqslant C_{4}\left\|\nabla S_{1}-\nabla S_{2}\right\|_{\infty}\|\varphi\|_{\infty}
$$

Proof. From $R_{S} L_{S}=\Pi$,

$$
\left(R_{S_{1}}-R_{S_{2}}\right) L_{S_{1}}+R_{S_{2}}\left(L_{S_{1}}-L_{S_{2}}\right)=0 .
$$

Multiplying this equality by $R_{S_{1}}$ on the right, and using that $R_{S} \Pi=R_{S}$, we get for all $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$,

$$
\left(R_{S_{1}}-R_{S_{2}}\right) \varphi=R_{S_{2}}\left(L_{S_{2}}-L_{S_{1}}\right) R_{S_{1}} \varphi .
$$

Conclusion follows from the estimate (27) and the fact that $\left\|e^{\beta S_{1}}-e^{\beta S_{2}}\right\|_{\infty} \leqslant C\left\|\nabla S_{1}-\nabla S_{2}\right\|_{\infty}$ for some $C>0$ uniformly over $S_{1}, S_{2} \in \mathcal{B}_{C_{0}}$.

Lemma 12. There exists $C_{5}>0$ such that for all $k \geqslant 1$ and $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$

$$
\left\|R_{A_{t_{k}}} \varphi-R_{A_{t_{k-1}}} \varphi\right\|_{\infty} \leqslant \frac{C_{5}}{k}\|\varphi\|_{\infty}
$$

Proof. From Lemma 5, $A_{t} \in \mathcal{B}_{C_{0}}$ for all $t \geqslant 0$, so that Lemma 11 applies. It remains to obtain a bound on $\left\|\nabla A_{t_{k}}-\nabla A_{t_{k-1}}\right\|_{\infty}$. In this proof, to simplify the notation, we write $\theta_{k}=\theta_{\nu_{t_{k}}}$ and $F_{k}=F_{\nu_{t_{k}}}$. Denoting by

$$
m:=\frac{\int_{t_{k}}^{t_{k+1}} \delta_{\left(Q_{s}, Z_{s}\right)} e^{\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s}{\int_{t_{k}}^{t_{k+1}} e^{\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s} \quad \text { and } \quad p:=\frac{\int_{t_{k}}^{t_{k+1}} e^{\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s}{\int_{0}^{t_{k+1}} e^{\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s},
$$

it holds that

$$
\nu_{t_{k+1}}=(1-p) \nu_{t_{k}}+p m .
$$

In particular, for some $c, c^{\prime}>0$, for all $z \in \mathbb{T}^{d}$ and $k \in \mathbb{N}$,

$$
\left|\theta_{k+1}(z)-\theta_{k}(z)\right|=\frac{p}{1+\lambda}\left|\int_{(q, y) \in \mathbb{T}^{p} \times \mathbb{T}^{d}} K(z, y)\left(\mathrm{d} m(q, y)-\mathrm{d} \nu_{t_{k}}(q, y)\right)\right| \leqslant c p \leqslant \frac{c^{\prime}}{k},
$$

where we used that $A_{s} \in \mathcal{B}_{C_{0}}$ for all $s \in\left[0, t_{k+1}\right]$. The same argument also works for the derivatives of $\theta_{\nu_{t}}$, for $F_{\nu_{t}}$ and its derivatives, so that for any multi-index $\alpha \in \mathbb{N}^{d}$, there exists a constant $C_{\alpha}$ such that for all $k \geqslant 1$,

$$
\left\|\partial^{\alpha} F_{k+1}-\partial^{\alpha} F_{k}\right\|_{\infty}+\left\|\partial^{\alpha} \theta_{k+1}-\partial^{\alpha} \theta_{k}\right\|_{\infty} \leqslant \frac{C_{\alpha}}{k}
$$

Now, from the Euler equations satisfied by $A_{t_{k}}$ and $A_{t_{k+1}}$, we get

$$
\begin{equation*}
\nabla \cdot\left(\theta_{k} \nabla\left(A_{t_{k}}-A_{t_{k+1}}\right)\right)=\nabla \cdot\left(\nabla A_{t_{k+1}}\left(\theta_{k+1}-\theta_{k}\right)\right)-\nabla \cdot\left(\theta_{k+1} F_{k+1}-\theta_{k} F_{k}\right) . \tag{28}
\end{equation*}
$$

Multiplying this equation by $A_{t_{k}}-A_{t_{k+1}}$, integrating and using Lemma 5 and the lower bound on $\theta_{k}$, we get

$$
\int_{\mathbb{T}^{d}}\left|\nabla\left(A_{t_{k}}-A_{t_{k+1}}\right)(z)\right|^{2} \mathrm{~d} z \leqslant \frac{c}{k^{2}}
$$

for some $c>0$. Next, differentiating (28), multiplying it by derivatives of $A_{t_{k}}-A_{t_{k+1}}$, integrating and using by induction the previous estimates, we obtain in fact that

$$
\int_{\mathbb{T}^{d}}\left|\nabla \partial^{\alpha}\left(A_{t_{k}}-A_{t_{k+1}}\right)(z)\right|^{2} \mathrm{~d} z \leqslant \frac{c_{\alpha}}{k^{2}}
$$

for some $c_{\alpha}>0$ for all $\alpha \in \mathbb{N}^{d}$, and Sobolev embeddings then yield the conclusion.

### 2.3 Proof of Proposition 7

The arguments are similar to those of the proof of [6, Theorem 3.6]. To alleviate notations, we write $R_{t}=R_{A_{t}}$ and, fixing some $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$, we set

$$
G_{t}(x)=\frac{1}{t} R_{t} \varphi(x),
$$

for all $x \in \mathbb{T}^{D}$. Recall the time-changed process $\left(\bar{X}_{t}\right)_{t \geqslant 0}$ solves (22). Denoting by $\tau_{k}=\tau\left(t_{k}\right)$, Itô's formula applied on the intervals $\left[\tau_{n}, \tau_{n+1}\right.$ [ yields, for $0<s<t$,

$$
\begin{aligned}
G_{t}\left(\bar{X}_{t}\right)-G_{s}\left(\bar{X}_{s}\right)= & M_{t}-M_{s}+\sum_{k \in \mathbb{N}} \mathbb{1}_{\left.\tau_{k} \in \mathrm{~J} s, t\right]} \frac{1}{\tau_{k}}\left(R_{\tau_{k}}-R_{\tau_{k-1}}\right) \varphi\left(\bar{X}_{\tau_{k}}\right) \\
& +\int_{s}^{t}\left(L_{u} G_{u}\left(\bar{X}_{u}\right)-\frac{1}{u^{2}} R_{u} \varphi\left(\bar{X}_{u}\right)\right) \mathrm{d} u
\end{aligned}
$$

where $M_{t}$ is a martingale with quadratic variation $2 \beta^{-1} \int_{s}^{t}\left|\nabla G_{u}\left(\bar{X}_{u}\right)\right|^{2} e^{\beta S_{u}\left(\bar{X}_{u}\right)} \mathrm{d} u$.
By definition of $R_{u}$,

$$
\int_{s}^{t} L_{u} G_{u}\left(\bar{X}_{u}\right) \mathrm{d} u=\int_{s}^{t} \frac{\varphi\left(\bar{X}_{u}\right)-\int \varphi \mathrm{d} \mu}{u} \mathrm{~d} u
$$

so that, recalling the definition (24) of $\varepsilon_{t}(s)$,

$$
\begin{aligned}
\left|\varepsilon_{t}(s) \varphi\right| & =\left|\int_{e^{t}}^{e^{t+s}} L_{u} G_{u}\left(\bar{X}_{u}\right) \mathrm{d} u\right| \\
& \leqslant\left|a_{1}(s, t)\right|+\left|a_{2}(s, t)\right|+\left|a_{3}(s, t)\right|+\left|a_{4}(s, t)\right|
\end{aligned}
$$

with

$$
\begin{aligned}
& a_{1}(s, t)=G_{e^{t+s}}\left(\bar{X}_{e^{t+s}}\right)-G_{e^{t}}\left(\bar{X}_{e^{t}}\right) \\
& a_{2}(s, t)=\int_{e^{t}}^{e^{t+s}} \frac{1}{u^{2}} R_{u} \varphi\left(\bar{X}_{u}\right) \mathrm{d} u \\
& a_{3}(s, t)=\sum_{k \in \mathbb{N}} \mathbb{1}_{\left.\left.\tau_{k} \in\right] e^{t}, e^{t+s}\right]} \frac{1}{\tau_{k}}\left(R_{\tau_{k}}-R_{\tau_{k-1}}\right) \varphi\left(\bar{X}_{\tau_{k}}\right) \\
& a_{4}(s, t)=M_{e^{t+s}}-M_{e^{t}}
\end{aligned}
$$

Lemma 5 states that $A_{t} \in \mathcal{B}_{C_{0}}$ for all $t \geqslant 0$ so that, for all $s, t \geqslant 0$, from (27),

$$
\begin{aligned}
& \left|a_{1}(s, t)\right| \leqslant 2 C_{3} e^{-t}\|\varphi\|_{\infty} \\
& \left|a_{2}(s, t)\right| \leqslant C_{3} e^{-t}\|\varphi\|_{\infty},
\end{aligned}
$$

and from Lemma 12 ,

$$
\left|a_{3}(s, t)\right| \leqslant C_{5}\|\varphi\|_{\infty} \sum_{k \in \mathbb{N}} \mathbb{1}_{\tau_{k}>e^{t}} \frac{1}{k \tau_{k}} .
$$

From the uniform bound on $A_{t}, \tau_{k}$ grows at least linearly, in the sense there exists $r_{1}>0$ such that, for all $k \in \mathbb{N}, \tau_{k+1}-\tau_{k}>r_{1}$. On the other hand, for the same reason, $\tau_{k}$ grows at most linearly, so that $\tau_{k} \leqslant r_{2} k$ for some $r_{2}$. Hence,

$$
\left|a_{3}(s, t)\right| \leqslant C_{5}\|\varphi\|_{\infty} \sum_{k \in \mathbb{N}} \frac{r_{2}}{\left(e^{t}+r_{1} k\right)^{2}} \leqslant c\|\varphi\|_{\infty} e^{-t}
$$

for some $c>0$.
Therefore, for all $T>0$, almost surely,

$$
\sup _{s \in[0, T]}\left(\left|a_{1}(s, t)\right|+\left|a_{2}(s, t)\right|+\left|a_{3}(s, t)\right|\right) \underset{t \rightarrow \infty}{\longrightarrow} 0
$$

The quadratic variation of $M_{e^{t+s}}-M_{e^{t}}$ being bounded by $2 \beta^{-1}\left(e^{\beta C_{0}} C_{3}\|\varphi\|_{\infty}\right)^{2} e^{-t}$, Doob's inequality implies that for all $T, \delta>0$,

$$
\mathbb{P}\left(\sup _{s \in[0, T]}\left|a_{4}(s, t)\right| \geqslant \delta\right) \leqslant \frac{2}{\beta \delta^{2}}\left(e^{\beta C_{0}} C_{3}\|\varphi\|_{\infty}\right)^{2} e^{-t}
$$

This implies

$$
\mathbb{P}\left(\sup _{s \in[0, T]}\left|a_{4}(s, t)\right| \geqslant e^{-\frac{t}{4}}\right) \leqslant \frac{2}{\beta}\left(e^{\beta C_{0}} C_{3}\|\varphi\|_{\infty}\right)^{2} e^{-\frac{t}{2}}
$$

and, by the Borel-Cantelli Theorem, almost surely,

$$
\limsup _{t \rightarrow \infty} \frac{1}{t} \log \sup _{s \in[0, T]}\left|a_{4}(s, t)\right| \leqslant-\frac{1}{4} .
$$

At the end of the day, we have proved that, for all $\varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right)$ and all $T>0$,

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \sup _{s \in[0, T]}\left|\varepsilon_{t}(s) \varphi\right|=0, \quad \text { almost surely. } \tag{29}
\end{equation*}
$$

Actually, in order to show that $\mathbb{R}_{+} \ni t \mapsto \bar{\nu}_{t}$ is an asymptotic pseudo-trajectory of the flow (25), the result we wish to prove is that

$$
\begin{equation*}
\text { almost surely, } \quad\left(\forall \varphi \in \mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right), \forall T>0, \quad \lim _{t \rightarrow \infty} \sup _{s \in[0, T]}\left|\varepsilon_{t}(s) \varphi\right|=0\right) \tag{30}
\end{equation*}
$$

Nevertheless, as stated in [6, Proposition 3.5], to get (30), it is sufficient to prove that (29) holds only for $\varphi$ and $T$ in a countable set of $\mathcal{C}^{\infty}\left(\mathbb{T}^{D}\right) \times \mathbb{R}_{+}$. This concludes the proof of Proposition 7, hence of Proposition 6, hence of Theorem 2.

## 3 Consistency of the tensor approximation

This section is devoted to the proof of Propositions 1 and 3 and Theorem 4. In all this section, Assumption 11 holds and we write $\theta=\theta_{\nu}, F=F_{\nu}, \mathcal{J}=\mathcal{J}_{\nu}$ for some fixed $\nu \in$ $\mathcal{P}\left(\mathbb{T}^{p} \times \mathbb{T}^{d}\right)$. Recall that $F, \theta \in \mathcal{C}^{\infty}\left(\mathbb{T}^{d}\right)$, that $\theta$ is a positive probability density on $\mathbb{T}^{d}$, and that the minimizers of $\mathcal{J}$ in $H$ are exactly the minimizers of $\widetilde{\mathcal{J}}$ in $H$, where for all $f \in H$,

$$
\widetilde{\mathcal{J}}(f)=\int_{\mathbb{T}^{d}}|F(z)-\nabla f(z)|^{2} \theta(z) \mathrm{d} z,
$$

the link between $\mathcal{J}$ and $\widetilde{\mathcal{J}}$ begin given by (11).
Since $\theta$ is bounded from above and below by positive constants, the weighted spaces $L^{2}\left(\mathbb{T}^{d} ; \theta\right)$ and $H^{1}\left(\mathbb{T}^{d} ; \theta\right)$ are equal to the flat spaces $L^{2}\left(\mathbb{T}^{d} ; \mathrm{d} z\right)$ and $H^{1}\left(\mathbb{T}^{d} ; \mathrm{d} z\right)$. We endow $H$ (whose definition is given in (4), with the norm

$$
\|f\|=\sqrt{(1+\lambda) \int_{\mathbb{T}^{d}}|\nabla f(z)|^{2} \theta(z) \mathrm{d} z}
$$

which is indeed a norm, equivalent to the usual $H^{1}$ norm from the Poincaré-Wirtinger inequality: there exists $C>0$ such that for all $f \in H$,

$$
\begin{aligned}
\int_{\mathbb{T}^{d}} f^{2}(z) \theta(z) \mathrm{d} z & \leqslant\|\theta\|_{\infty} \int_{\mathbb{T}^{d}} f^{2}(z) \mathrm{d} z \\
& \leqslant C\|\theta\|_{\infty} \int_{\mathbb{T}^{d}}|\nabla f(z)|^{2} \mathrm{~d} z \\
& \leqslant C\|\theta\|_{\infty}\left\|\theta^{-1}\right\|_{\infty} \int_{\mathbb{T}^{d}}|\nabla f(z)|^{2} \theta(z) \mathrm{d} z
\end{aligned}
$$

The scalar product associated with $\|\cdot\|$ is denoted by $\langle\cdot\rangle$. The choice of such a norm is motivated by the fact that, denoting by $\mathcal{J}^{\prime}$ the differential of $\mathcal{J}$, then for all $f, g \in H$,

$$
\begin{aligned}
\mathcal{J}(f) & =\mathcal{J}(0)+\mathcal{J}^{\prime}(0) \cdot f+\|f\|^{2} \\
\mathcal{J}^{\prime}(f) \cdot g & =\mathcal{J}^{\prime}(0) \cdot g+2\langle f, g\rangle .
\end{aligned}
$$

Proposition 1 is then a direct consequence of the strict convexity of $\mathcal{J}$. The unique minimizer $f_{*}$ of $\mathcal{J}$ over $H$ satisfies

$$
\forall g \in H, \quad \int F(z) \cdot \nabla g(z) \theta(z) \mathrm{d} z=\int \nabla f_{*}(z) \cdot \nabla g(z) \theta(z) \mathrm{d} z
$$

or equivalently

$$
\begin{equation*}
\forall g \in H, \quad \mathcal{J}(g)=\mathcal{J}\left(f_{*}\right)+\left\|f_{*}-g\right\|^{2} . \tag{31}
\end{equation*}
$$

Proof of Proposition 3. Let $f \in H$ and $i \in \llbracket 1, d \rrbracket$. If $\mathcal{J}(f)=\inf \left\{\mathcal{J}(f+g), g \in \Sigma_{i}\right\}$ then the result is correct since $0 \in \Sigma_{i}$ is a minimizer over $\Sigma_{i}$. Suppose now that 0 is not a minimizer, i.e. that $\mathcal{J}(f)>\inf \left\{\mathcal{J}(f+g), g \in \Sigma_{i}\right\}$ and consider a minimizing sequence $\left(g^{(l)}\right)_{l \in \mathbb{N}}$ in $\Sigma_{i}$ such that $\mathcal{J}\left(f+g^{(l)}\right)$ converges to $\inf \left\{\mathcal{J}(f+g), g \in \Sigma_{i}\right\}$ as $l$ goes to infinity. For $l$ large enough, $\mathcal{J}\left(f+g^{(l)}\right)<\mathcal{J}(f)$ so that $g^{(l)} \neq 0$, and thus up to an extraction we suppose that $g^{(l)} \neq 0$ for all $l \in \mathbb{N}$. Moreover the sequence is bounded in $H^{1}$ and thus, up to the extraction of a subsequence, we suppose that it weakly converges in $H^{1}$ to some $g^{*} \in H$. The function $H \ni g \mapsto \mathcal{J}(f+g)$ being convex on $H$, it is weakly lower semi-continuous, so that

$$
\mathcal{J}\left(g_{*}\right) \leqslant \inf _{g \in \Sigma_{i}} \mathcal{J}(f+g)
$$

For all $l \in \mathbb{N}$, there exist $r_{1}^{(l)}, \cdots, r_{d}^{(l)} \in H^{1}(\mathbb{T})$ such that $g^{(l)}=\bigotimes_{j=1}^{d} r_{j}^{(l)}$. Since $g^{(l)} \neq 0$, we can normalize the $r_{j}^{(l)}$ s so that $\left\|r_{j}^{(l)}\right\|_{L^{2}(\mathbb{T})}=1$ for all $j \neq i$ and $l \in \mathbb{N}$. As a consequence, up to the extraction of a subsequence, for all $1 \leq j \neq i \leq d$, there exists $r_{j}^{*} \in L^{2}(\mathbb{T})$ such that the sequence $\left(r_{j}^{(l)}\right)_{l \in \mathbb{N}}$ weakly converges to $r_{j}^{*}$ in $L^{2}(\mathbb{T})$. Now, since the sequence $g^{(l)}$ is bounded in $H$ and

$$
\left\|\nabla g^{(l)}\right\|_{L^{2}\left(\mathbb{T}^{d}\right)}^{2}=\sum_{j=1}^{d}\left\|\partial_{z_{j}} r_{j}^{(l)}\right\|_{L^{2}(\mathbb{T})}^{2} \prod_{h \neq j}\left\|r_{h}^{(l)}\right\|_{L^{2}(\mathbb{T})}^{2},
$$

we get that the sequence $\left(r_{i}^{(l)}\right)_{l \in \mathbb{N}}$ is bounded in $H^{1}(\mathbb{T})\left(\right.$ since $\int_{\mathbb{T}} r_{i}^{(l)}=0$ for all $\left.l \in \mathbb{N}\right)$. Thus, up to the extraction of another subsequence, the sequence $\left(r_{i}^{(l)}\right)_{l \in \mathbb{N}}$ weakly converge in $H^{1}(\mathbb{T})$ to some $r_{i}^{*} \in H^{1}(\mathbb{T})$ such that $\int_{\mathbb{T}} r_{i}^{*}=0$. From [17, Lemma 2], $\left(g^{(l)}\right)_{l \in \mathbb{N}}$ converges in the distributional sense to $\bigotimes_{j=1}^{d} r_{j}^{*}$. Thus, $g^{*}=\bigotimes_{j=1}^{d} r_{j}^{*}$ and since $g^{*} \neq 0$, this implies that for all $1 \leq j \leq d, r_{j}^{*} \neq 0$. Finally, since

$$
\left\|\nabla g^{*}\right\|_{L^{2}\left(\mathbb{T}^{d}\right)}^{2}=\sum_{j=1}^{d}\left\|\partial_{z_{j}} r_{j}^{*}\right\|_{L^{2}(\mathbb{T})}^{2} \prod_{h \neq j}\left\|r_{h}^{*}\right\|_{L^{2}(\mathbb{T})}^{2}
$$

is a finite quantity, this implies that for all $1 \leq j \neq i \leq d$,

$$
\left\|\partial_{z_{j}} r_{j}^{*}\right\|_{L^{2}(\mathbb{T})} \leq \frac{\left\|\nabla g^{*}\right\|_{L^{2}\left(\mathbb{T}^{d}\right)}^{2}}{\prod_{h \neq j}\left\|r_{h}^{*}\right\|_{L^{2}(\mathbb{T})}^{2}}<+\infty
$$

and thus $r_{j}^{*} \in H^{1}(\mathbb{T})$. This implies that $g_{*} \in \Sigma_{i}$ and yields the desired result.
Remark 13. The problem would be ill-posed if we were to try and minimize $\mathcal{J}\left(f+g-\int_{\mathbb{T}^{d}} g\right)$ over all $g \in \Sigma:=\left\{r_{1} \otimes \cdots \otimes r_{d}, r_{j} \in H^{1}(\mathbb{T})\right.$ for all $\left.1 \leq j \leq d\right\}$. This is the reason why we introduced the condition that one of the $r_{j}$ 's has zero mean. Indeed, consider the situation where $d=2, f=0$ and $F(z)=\left(a^{\prime}\left(z_{1}\right), b^{\prime}\left(z_{2}\right)\right)$ for some smooth functions $a, b: \mathbb{T} \rightarrow \mathbb{R}$ with zero
mean. Then, the minimum of $\mathcal{J}$ over $H$ is 0 , and only attained at $f^{*}\left(z_{1}, z_{2}\right)=a\left(z_{1}\right)+b\left(z_{2}\right)$, which is not of the form $g-\int_{\mathbb{T}^{d}} g$ for some $g \in \Sigma$. Nevertheless, the sequence $\left(g^{(l)}\right)_{l \in \mathbb{N}^{*}}$ defined by: for all $l \in \mathbb{N}^{*}, g^{(l)}=r_{1}^{(l)} r_{2}^{(l)}$ with

$$
r_{1}^{(l)}\left(z_{1}\right)=1+\frac{a\left(z_{1}\right)}{l}, \quad r_{2}^{(l)}\left(z_{2}\right)=l+b\left(z_{2}\right)
$$

is a minimizing sequence. Indeed, the sequence $\left(g^{(l)}-\int_{\mathbb{T}^{d}} g^{(l)}\right)_{l \in \mathbb{N}^{*}}$ weakly converges to $f^{*}$. In other words, the set $\left\{g-\int_{\mathbb{T}^{d}} g, g \in \Sigma\right\}$ is not weakly closed in $H$.

Proposition 3 proves that all the iterations of Algorithm 2 are well-defined. In the following, we consider a sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ given by the latter and $g_{n}=f_{n+1}-f_{n}$ for $n \in \mathbb{N}$. The general idea of the proof of Theorem 4 is that, if the sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ converges to some $f_{\infty}$ in $H$, it holds that $\mathcal{J}^{\prime}\left(f_{\infty}\right) \cdot g=0$ for all $g \in \cup_{i=1}^{d} \Sigma_{i}$, and a density argument enables to conclude. Nevertheless, remark that $\Sigma_{i}$ is not a vector space and that its elements all have null integral, so that one should be careful. In the following, we essentially adapt the arguments of 9].
Lemma 14. For all $g \in \cup_{i=1}^{d} \Sigma_{i}$ and $n \in \mathbb{N}$,

$$
\left|\mathcal{J}^{\prime}\left(f_{n}\right) \cdot g\right| \leqslant 6\|g\| \sum_{j=0}^{d-1}\left\|g_{n+j}\right\|
$$

Proof. Let $i \in \llbracket 1, d \rrbracket$ be such that $g \in \Sigma_{i}$. Let $n_{i} \in \llbracket n, n+d-1 \rrbracket$ be such that $g_{n_{i}} \in \Sigma_{i}$. We bound, first,

$$
\left|\mathcal{J}^{\prime}\left(f_{n}\right) \cdot g\right| \leqslant\left|\mathcal{J}^{\prime}\left(f_{n_{i}}\right) \cdot g\right|+2\left|\left\langle g, f_{n_{i}}-f_{n}\right\rangle\right| .
$$

The second term of the right hand side is bounded by $2\|g\| \sum_{j=n}^{n_{i}-1}\left\|g_{j}\right\|$. To deal with the first one, note that, even though $g_{n_{i}}$ is a minimizer of $\mathcal{J}\left(f_{n_{i}}+\cdot\right)$ over $\Sigma_{i}$, it is not necessarily true that $\mathcal{J}^{\prime}\left(f_{n_{i}}+g_{n_{i}}\right) \cdot g=0$, since $\Sigma_{i}$ is not a vector space. We follow the proof of [9, Proposition 3.3]. By convexity of

$$
t \in \mathbb{R} \mapsto \psi(t):=\mathcal{J}\left(f_{n_{i}}+g+t\left(g_{n_{i}}-g\right)\right),
$$

and since $t=0$ minimizes $\psi(t)$, we get

$$
\psi^{\prime}(0) \leqslant \psi(1)-\psi(0) \leqslant 0
$$

which reads

$$
\mathcal{J}^{\prime}\left(f_{n_{i}}+g\right) \cdot g \geqslant \mathcal{J}^{\prime}\left(f_{n_{i}}+g\right) \cdot g_{n_{i}} .
$$

Hence,

$$
\begin{aligned}
-\mathcal{J}^{\prime}\left(f_{n_{i}}\right) \cdot g & =-\mathcal{J}^{\prime}\left(f_{n_{i}}+g\right) \cdot g+2\|g\|^{2} \\
& \leqslant-\mathcal{J}^{\prime}\left(f_{n_{i}}+g\right) \cdot g_{n_{i}}+2\|g\|^{2} \\
& \leqslant-\mathcal{J}^{\prime}\left(f_{n_{i}}+g_{n_{i}}\right) \cdot g_{n_{i}}+2\left\langle g_{n_{i}}, g_{n_{i}}-g\right\rangle+2\|g\|^{2} .
\end{aligned}
$$

Now, 1 being a minimizer over $\mathbb{R}$ of $t \mapsto \mathcal{J}\left(f_{n_{i}}+t g_{n_{i}}\right), \mathcal{J}^{\prime}\left(f_{n_{i}}+g_{n_{i}}\right) \cdot g_{n_{i}}=0$, so that

$$
-\mathcal{J}^{\prime}\left(f_{n_{i}}\right) \cdot g \leqslant 2\left(\left\|g_{n_{i}}\right\|^{2}+\|g\|\left\|g_{n_{i}}\right\|+\|g\|^{2}\right) .
$$

When applied to $\widetilde{g}= \pm\left\|g_{n_{i}}\right\| g /\|g\|$, this inequality yields

$$
\left|\mathcal{J}^{\prime}\left(f_{n_{i}}\right) \cdot g\right| \leqslant 6\|g\|\left\|g_{n_{i}}\right\|,
$$

which concludes the proof.

Proposition 15. Let $f_{*}$ be the unique minimizer of $\mathcal{J}$ over $H$. Then

$$
\left\|f_{n}-f_{*}\right\| \underset{n \rightarrow \infty}{\longrightarrow} 0
$$

Proof. As in the previous proof,

$$
0=\mathcal{J}^{\prime}\left(f_{n}+g_{n}\right) \cdot g_{n}=2(1+\lambda) \int \nabla g_{n} \cdot\left(\nabla f_{n+1}-F\right) \theta=0
$$

for all $n \in \mathbb{N}$, so that,

$$
\begin{equation*}
\mathcal{J}\left(f_{n}\right)-\mathcal{J}\left(f_{n+1}\right)=\mathcal{J}\left(f_{n+1}-g_{n}\right)-\mathcal{J}\left(f_{n+1}\right)=\left\|g_{n}\right\|^{2} \tag{32}
\end{equation*}
$$

In particular, since $\left(\mathcal{J}\left(f_{n}\right)\right)_{n \geqslant 0}$ is a decreasing sequence bounded from below,

$$
\begin{equation*}
\sum_{n \geqslant 0}\left\|g_{n}\right\|^{2}<\infty \tag{33}
\end{equation*}
$$

Together with Lemma 14 and the fact $\mathcal{J}^{\prime}\left(f_{*}\right)=0$, this implies that for all $g \in \operatorname{Span}\left(\cup_{i=1}^{d} \Sigma_{i}\right)$,

$$
2\left\langle f_{*}-f_{n}, g\right\rangle=\mathcal{J}^{\prime}\left(f_{*}\right) \cdot g-\mathcal{J}^{\prime}\left(f_{n}\right) \cdot g \underset{n \rightarrow \infty}{\longrightarrow} 0
$$

Now, for all $r_{1}, \cdots, r_{d} \in \mathcal{C}^{\infty}(\mathbb{T})$, denoting by $h:=\bigotimes r_{i} \in \mathcal{C}^{\infty}\left(\mathbb{T}^{d}\right)$ (note that we do not have necessarily that $\int_{\mathbb{T}^{d}} h=0$ ), we can write

$$
h-\int_{\mathbb{T}^{d}} h=\sum_{j=1}^{d}\left(r_{j}-\int_{\mathbb{T}} r_{j}\right)\left(\prod_{l<j} \int_{\mathbb{T}} r_{l}\right) \prod_{l>j} r_{l}
$$

which proves that $h-\int_{\mathbb{T}^{d}} h \in \operatorname{Span}\left(\cup_{i=1}^{d} \Sigma_{i}\right)$. As a consequence,

$$
(1+\lambda) \int_{\mathbb{T}^{d}} \theta \nabla\left(f_{*}-f_{n}\right) \cdot \nabla h=\left\langle f_{*}-f_{n}, h-\int_{\mathbb{T}^{d}} h\right\rangle \underset{n \rightarrow \infty}{\longrightarrow} 0,
$$

As a consequence, the limit $f_{\infty}$ of any convergent (in the weak sense in $H$ ) subsequence of $\left(f_{n}\right)_{n \in \mathbb{N}}$ necessarily satisfies that

$$
\int_{\mathbb{T}^{d}} \theta \nabla\left(f_{*}-f_{\infty}\right) \cdot \nabla h=0
$$

for any tensor product function $h=\bigotimes r_{i}$, with $r_{i} \in \mathcal{C}^{\infty}(\mathbb{T})$ for all $1 \leq i \leq d$. By [9, Lemma 2.1], this implies that $f_{\infty}=f_{*}$. On the other hand, since $\left(\left\|f_{n}\right\|\right)_{n \in \mathbb{N}}$ is bounded, all its subsequences admits weak convergent subsequences, and the fact they all have the same limit $f_{*}$ proves that the whole sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ weakly converges in $H$ to $f_{*}$. In particular

$$
\begin{equation*}
\left\langle f_{*}-f_{n}, f_{*}\right\rangle \underset{n \rightarrow \infty}{\longrightarrow} 0 \tag{34}
\end{equation*}
$$

Thus, it only remains to prove that $\left(\left\langle f_{*}-f_{n}, f_{n}\right\rangle\right)_{n \in \mathbb{N}}$ also converges to zero as $n$ goes to infinity to obtain the strong convergence of the sequence $\left(f_{n}\right)_{n \in \mathbb{N}}$ to $f_{*}$ in $H$. From Lemma 14 ,

$$
\begin{aligned}
2\left|\left\langle f_{*}-f_{n}, f_{n}\right\rangle\right| & =\left|\mathcal{J}^{\prime}\left(f_{n}\right) \cdot f_{n}\right| \\
& \leqslant \sum_{k=0}^{n}\left|\mathcal{J}^{\prime}\left(f_{n}\right) \cdot g_{k}\right| \\
& \leqslant 6\left(\sum_{k=0}^{n-1}\left\|g_{k}\right\|\right) \sum_{j=n}^{n+d-1}\left\|g_{j}\right\| \\
& \leqslant 6 \sqrt{n d a_{n} \sum_{k=0}^{\infty}\left\|g_{k}\right\|^{2} .}
\end{aligned}
$$

with $a_{n}=\sum_{j=n}^{n+d-1}\left\|g_{j}\right\|^{2}$. Using $\sqrt{33}$, since $\sum_{n \in \mathbb{N}} a_{n} \leq d \sum_{n \in \mathbb{N}}\left\|g_{n}\right\|^{2}<\infty$, there exists an extracted subsequence $\left(n_{k}\right)_{k \geqslant 1}$ such that $\left(n_{k} a_{n_{k}}\right)_{k \geq 1}$ converges to 0 as $k$ goes to infinity. As a consequence, $\left(\left\langle f_{*}-f_{n_{k}}, f_{n_{k}}\right\rangle\right)_{k \geq 1}$ goes to zero as $k \rightarrow \infty$, and thus so does $\left(\left\|f_{n_{k}}-f_{*}\right\|\right)_{k \geq 1}$ by (34). Finally, from (31), the sequence $\left(\left\|f_{n}-f_{*}\right\|\right)_{n \in \mathbb{N}}$ is non-increasing, so that the whole sequence goes to zero. Hence the result.

## 4 Discussion and variations

For the sake of clarity, the TABF algorithm defined in Section 1.3 has been kept relatively simple, and there is obviously room for many variations or fine-tuning. We list here a few of them.

### 4.1 Extended ABF

Consider general reaction coordinates $\xi: \mathbb{T}^{D} \rightarrow \mathcal{M}$ where $\mathcal{M}$ is a submanifold of $\mathbb{R}^{d}$ or $\mathbb{T}^{d}$. In the Extended ABF (EABF) algorithm introduced in [13] (see also [21]), the state space is extended to $\mathbb{T}^{D} \times \mathcal{M}$ with the addition of auxiliary variables (or fictitious particles) $z \in \mathcal{M}$, and the potential $V$ on $\mathbb{T}^{D}$ is extended to a potential $\tilde{V}$ on $\mathbb{T}^{D} \times \mathcal{M}$ as

$$
\tilde{V}(q, z)=V(x)+\frac{1}{2 \sigma^{2}}\left(\operatorname{dist}_{\mathcal{M}}(\xi(q), z)\right)^{2}, \quad \forall(q, z) \in \mathbb{T}^{D} \times \mathcal{M}
$$

for some small parameter $\sigma>0$, where dist ${ }_{\mathcal{M}}$ stands for the distance on $\mathcal{M}$. The reaction coordinates on the extended space are then defined by $\tilde{\xi}(q, z)=z$, which means the framework considered in the present paper is general for the EABF algorithm. If $(Q, Z)$ is distributed according to $\mu_{\tilde{V}, \beta}$, the law of $Z$ is obtained from the law of $\xi(Q)$ through a Gaussian convolution of variance $\sigma^{2} / \beta$. There are several practical advantages to EABF:

- In the potential $\tilde{V}$, in the case where dist ${ }_{\mathcal{M}}$ is an Euclidian distance, the $z_{i}$ 's for $i \in \llbracket 1, d \rrbracket$ appear in separate terms of a sum, they are not directly coupled. As a consequence, in the EABF case, $\mu_{A, \beta}$ should be, in some sense, closer to the product of its marginal (namely, at equilibrium, the $Z_{i}$ 's should be closer to be independent) than in the non-extended ABF case. In [13], this was a crucial point since the density $\mu_{A, \beta}$ was approximated by a tensor product. But even in our case where the approximation as a sum of tensor product is made at the level of $A$, we can expect this form of $\tilde{V}$ to improve the approximation.
- After convolution, the so-called mean-force $\nabla_{z} A$ is smoother than the initial mean force in the non-extended ABF case. Since it varies less, its estimation is expected to be easier.

That being said, the tensorized ABF introduced above can also be straightforwardly extended to a general ABF framework, without extended coordinates.

### 4.2 Non-periodic reaction coordinates.

In general, $\mathcal{M}$ may be different from $\mathbb{T}^{d}$. If it has boundaries, for instance if $\mathcal{M}=[0,1]^{d}$, the definition of the algorithm is the same except that the diffusion (3) is reflected at the boundaries of $\mathcal{M}$. The proof of well-posedness and convergence of the tensor algorithm, i.e. Theorem 4, is unchanged. The proof of the long-time convergence of the idealized algorithm, i.e. Theorem 2, is similar up to technical considerations in particular to take into account boundary conditions in Section 2.2.

Moreover, $\mathcal{M}$ may not be compact, for instance $\mathcal{M}=\mathbb{R}^{d}$. Since the Lebesgue measure has not a finite mass, a confining biasing potential has to be added to the adaptive biasing potential, see [19, Section 1.2], in which case the law of $\xi\left(X_{t}\right)$ does not converge to a uniform law (flat histogram) but to a target unimodal law on $\mathbb{R}^{d}$.

### 4.3 Real implementation.

The algorithm really implemented for the numerical experiments in Section 5 differs from the theoretical Algorithm 3 in the following points:

1. Time and space are discretized. The SDE (3) is replaced by an Euler-Maruyama scheme with some timestep $\delta t$ and the time integral in (6) is replaced by a discrete sum with a timestep $\Delta t$ (not necessarily small; it can be of the order of the decorrelation length of the process $\left.\left(Q_{t}, Z_{t}\right)_{t \geqslant 0}\right)$. The one-dimensional functions in the tensor terms are restricted to be continuous piecewise linear, determined by their value on a discrete grid with $q \in \mathbb{N}_{*}$ points, so that solving (17) amount to solve a $q \times q$ linear system. In particular, the discrete space interpolation plays a role similar to the regularization kernel $K$ which is no more necessary, hence is discarded.
2. In fact, it is not necessary to re-weight the occupation distribution, namely (6) can be replaced by $\nu_{t}=1 / t \int_{0}^{t} \delta_{\left(Q_{s}, Z_{s}\right)} \mathrm{d} s$. In that case, $\nu_{t}$ is expected to converge to $\mu_{V-\tilde{A}_{*}, \beta}$ for some $\tilde{A}_{*}$ instead of $\mu_{V, \beta}$ but the conditional law of $Q$ given $Z=z$ is the same for these two laws. Since the free energy only depends on these conditional laws, $A_{t}$ is still expected to converge to $\tilde{A}_{*}$, that should be close to the true free energy in a sense similar to (10). This is clear in the mean-field limit of the algorithm, where no regularization is needed so that $\tilde{A}_{*}=A$ (see [19]). It is more difficult to establish for the self-interacting ABF process, but in parallel of the present paper it has been done in [5]. We tried numerically both cases, and the results were similar. The results presented in Section 5 are obtained with the full (non reweighted) occupation distribution.
3. Instead of a single particle, in practice, several replicas of the process (3) are simulated in parallel. Denoting $N$ the number of replicas and $\left(Q_{t}^{i}, Z_{t}^{i}\right)_{t \geqslant 0}$ the $i^{\text {th }}$ replica, $i \in \llbracket 1, N \rrbracket$, the total empirical distribution of the system is

$$
\begin{equation*}
\tilde{\nu}_{N, t}=\frac{1}{N\lfloor t / \Delta t\rfloor} \sum_{i=1}^{N} \sum_{k=1}^{\lfloor t / \Delta t\rfloor} \delta_{\left(Q_{k \Delta t}^{i}, Z_{k \Delta t}^{i}\right)} . \tag{35}
\end{equation*}
$$

The replicas all use the same bias $A_{t}$ obtained from this empirical distribution by minimizing $\mathcal{J}_{\tilde{\nu}_{N, t}}$ at times $t=t_{k}=k T_{u p}$.

### 4.4 Other possible simple variations.

1. From the biased trajectory $\left(Q_{t}, Z_{t}\right)_{t \geqslant 0}$ provided by the TABF algorithm, in order to compute expectations with respect to the target Gibbs measure $\mu=\mu_{V, \beta}$, an alternative to the reweighting step (2) is the following. Remark that only the $Z$ variable is biased, so that for all $z \in \mathbb{T}^{d}$, the conditional expectations $\int_{\mathbb{T}^{p}} f(q, z) \mu(q, z) \mathrm{d} q / \int_{\mathbb{T}^{p}} \mu(q, z) \mathrm{d} q$ can be estimated without re-weighting. On the other hand, the marginal law of $Z$ is estimated by $\exp \left(-\beta A_{T_{\text {tot }}}\right) / \int_{\mathbb{T}^{d}} \exp \left(-\beta A_{T_{\text {tot }}}(z)\right) \mathrm{d} z$.
2. The bias update period $T_{u p}$ and the number $m$ of tensor products added at each update in Algorithm 3, instead of having fixed values, could be adaptively chosen. For instance,
the bias could be updated when the histogram of the reaction coordinates have reached some stability, and $m$ could be the lowest integer $n \in \mathbb{N}$ such that $\mathcal{J}_{\nu}\left(A_{t_{k}}+f_{n-d}\right)$ $\mathcal{J}_{\nu}\left(A_{t_{k}}+f_{n}\right) \leqslant \varepsilon$ for some threshold $\varepsilon>0$.
3. A time-dependent weight in the definition (6) of $\nu_{t}$ (or, in practice, in (35) for $\tilde{\nu}_{N, t}$ ) can be added in such a way that old samples have less influence than new ones since they are more biased toward the initial distribution.
4. The regularization kernel $K$ and parameter $\lambda$ may depend on time. Indeed, as time goes, the size of the sample increases. Since the problem of minimizing $\mathcal{J}_{\nu_{t}}$ is in practice solved on a finite dimension space, for a time large enough the regularization is actually not necessary anymore and the minimization problem with $K(z, y)=\delta_{z}(y)$ and $\lambda=0$ is well-posed.
5. It is possible to use the tensor approximation only as a correction of the classical ABF , or more precisely of the Generalized ABF (GABF) algorithm proposed in [27] where the bias is just a sum of one-dimensional functions. Namely, for all time $t>0$, for $j \in \llbracket 1, d \rrbracket$, let

$$
\begin{aligned}
\alpha_{j, t}\left(z_{j}\right) & =\int_{\mathbb{T}^{p} \times \mathbb{T}^{d}} \partial_{y_{k}} V(q, y) K\left(y_{j}, z_{j}\right) \mathrm{d} \nu_{t}(q, y) \\
\beta_{j, t}\left(z_{j}\right) & =\int_{\mathbb{T}^{p} \times \mathbb{T}^{d}} K\left(y_{j}, z_{j}\right) \mathrm{d} \nu_{t}(q, y) .
\end{aligned}
$$

These functions can be recorded on $d$ one-dimensional grids and are easily updated on the fly. Denoting $\gamma_{t, j}\left(z_{j}\right)=\mathbb{1}_{\beta_{j, t}\left(z_{j}\right)>s} \alpha_{j, t}\left(z_{j}\right) / \beta_{j, t}\left(z_{j}\right)$ for some burn-in time $s>0$, let $A_{t, j}\left(z_{j}\right)=\int_{0}^{z_{j}} \gamma_{t, j}(z) \mathrm{d} z$ if the $j^{\text {th }}$ reaction coordinate $z_{j}$ lies in $\mathbb{R}$ and $A_{t, j}\left(z_{j}\right)=$ $\int_{0}^{z_{j}} \gamma_{t, j}(u) \mathrm{d} u-z_{j} \int_{0}^{1} \gamma_{t, j}(u) \mathrm{d} u$ if $z_{j}$ lies in $\mathbb{T}$ (so that, in both cases, $\partial_{z_{j}} A_{t, j}$ is the Helmoltz projection in $L^{2}\left(\mathrm{~d} z_{j}\right)$ of $\left.\gamma_{t, j}\right)$. Then, at time $t$, in the dynamics (3), use the bias $\nabla_{z} A_{t}$ with $A_{t}(z)=\sum_{j=1}^{d} A_{t_{k}, j}\left(z_{j}\right)+f_{m}(z)$ where $f_{m}$ is a tensor approximation obtained through Algorithm 2 of the minimizer of $H \ni f \mapsto \mathcal{J}_{\nu_{t_{k}}}\left(f-\sum_{j=1}^{d} A_{t_{k}, j}\right)$, where $t_{k}=\sup \left\{t_{k^{\prime}}<\right.$ $\left.t, k^{\prime} \in \mathbb{N}\right\}$ is the last update time.
6. For $k \in \mathbb{N}_{+}$, denote by $\mathcal{J}_{k}^{\lambda}$ the function given by (9) for some $\lambda>0$ with $\nu=\nu_{t_{k}}$. Rather than setting $A_{t_{k}}$ to be the minimizer of $\mathcal{J}_{k}^{\lambda}$, we can set it to be $A_{t_{k-1}}+f$ where $f$ is the minimizer of

$$
\begin{equation*}
H \ni f \mapsto \mathcal{J}_{k}^{0}\left(A_{t_{k-1}}+f\right)+\lambda \int_{\mathbb{T}^{d}}\left|\nabla_{z} f(z)\right|^{2} \mathrm{~d} z \tag{36}
\end{equation*}
$$

the difference being that $A_{t_{k-1}}$ does not appear in the last regularization term any more. Note that, when $\lambda=0$, there is no difference. When $\lambda>0$, the theoretical results of Section 3, i.e. the well-posedness of the tensor approximation, can be straightforwardly adapted. The long-time behaviour study of Section 2 should be similar, although a bit more troublesome since $A_{t_{k}}$ would not depend only on the empirical distribution $\eta_{t_{k}}$ but also on the previous bias $A_{t_{k-1}}$. On the other hand, remark that 0 is a minimizer of (36) if and only if $A_{t_{k}}$ is a minimizer of $\mathcal{J}_{k}^{0}$. As a consequence, the long-time limit of $A_{t}$ should be the minimizer of $\mathcal{J}_{\mu}$ with $\lambda=0$ which, in view of (10), advocates for this alternative form of cost function.
7. Since the bias is stored in memory in a tensor form, it is possible to use at some times a compression algorithm (see [11]) to reduce the number of tensor terms, if needed.

### 4.5 Some limitations and perspectives.

A practical limitation observed in the algorithm is the following. Recall that $d$ is too large to keep in memory the empirical measure on a grid by simply recording how many times each $d$-dimensional cell has been visited by the process, as in the classical ABF algorithm. Instead, the sequence $\left(Z_{k \Delta t}, \nabla_{z} V\left(Q_{k \Delta t}, Z_{k \Delta t}\right)\right)_{k \in \mathbb{N}}$ is kept in memory for some $\Delta t>0$, and thus computing an expectation with respect to $\nu_{t}$ has a numerical cost proportional to $t$. Such integrals are computed when solving the one-dimension equations (17), which have to be solved repeatedly at each addition of a tensor term to the bias. As $t$ grows, the update of the bias gets numerically more expensive. We list here some possible directions to address this question. The analysis of these variations is beyond the reach of the present work.

1. At the beginning of Algorithm 2, a clustering or quantization algorithm (see [25]) can be used to reduce the memory $\left(Z_{k \Delta t}, \nabla_{z} V\left(Q_{k \Delta t}, Z_{k \Delta t}\right)\right)_{k \in \llbracket 1, t_{n} / \Delta t \rrbracket}$ to fewer points.
2. Another way to deal with this problem would be to use a fixed small size for the memory. For instance, at an update time $t_{k}$, the empirical measure used to define $\mathcal{J}_{\bar{\nu}_{t_{k}}}$ could be

$$
\bar{\nu}_{t_{k}}=\left(\int_{t_{k-l}}^{t_{k}} e^{-\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s\right)^{-1} \int_{t_{k-l}}^{t_{k}} \delta_{\left(Q_{s}, Z_{s}\right)} e^{-\beta A_{s}\left(Z_{s}\right)} \mathrm{d} s
$$

for some small $l \in \mathbb{N}_{*}$, say $l=1$. In that case, in order to expect a long-time convergence of the bias, following classical stochastic algorithms, we would define the new bias as $A_{t_{k}}=A_{t_{k-1}}+\gamma_{k} f_{k}$ where $f_{k}$ is (a tensor approximation of) a minimizer over $H$ of $H \ni f \mapsto \mathcal{J}_{\nu_{t_{k}}}\left(A_{t_{k}}+f\right)$ and $\left(\gamma_{k}\right)_{k \in \mathbb{N}}$ is a positive sequence with $\gamma_{k} \rightarrow 0$ and $\sum_{l=1}^{k} \gamma_{l} \rightarrow \infty$ as $k \rightarrow \infty$.
3. Finally, a third way to deal with the memory management as time increases could be to use a stochastic gradient descent when solving the one-dimensional partial differential equation (17). In other words, when optimizing $r_{i}$ for some $1 \leq i \leq d$, instead of computing averages over all steps $l \in \llbracket 1, t_{k} / \Delta t \rrbracket$, only use an approximation of $\nu_{t_{k}}$ by picking a random (and comparatively small) set of steps among $\llbracket 1, t_{k} / \Delta t \rrbracket$. Then only an estimation of the gradient of $H^{1}(\mathbb{T}) \ni r_{i} \mapsto \mathcal{J}_{t_{k}}\left(f+\bigotimes_{j=1}^{d} r_{j}\right)$ is computed, which is exactly the settings of the stochastic gradient descent.

A second possible limitation is the following. Note that, as the number of reaction coordinates increases, we can expect that, at some point, the biasing scheme becomes unefficient. Indeed, by flattening the energy landscape, we replace the initial sampling problem (that was mainly restricted to low-energy regions, which form a low-dimensional manifold) by the sampling of the uniform measure on some hypercube, which is not so easy. In some sense, following the definitions of [18], at some point, energy barriers are replaced by entropic ones. Moreover, the variance of the estimator (2) increases due to the exponential weights. As a consequence, as $d$ increases, a partial biasing with $V_{\text {bias }, t}=\theta A_{t} \circ \xi$ for some $\theta \in(0,1)$ may be more appropriate than the full biasing (i.e. $\theta=1$ ). At the biased equilibrium, if $A_{t}=A$ is the true free energy, the marginal law of the reactions coordinates is thus $\mu_{(1-\theta) A, \beta}$, i.e. the temperature is increased. Then the choice of $\theta$ such that this measure satisfies a Poincaré inequality with minimal constant (which means the corresponding overdamped Langevin process mixes the fastest) may not be $\theta=1$.

## 5 Numerical experiments

Let us fix some details and parameters that will hold for the different examples below. In this section, the modifications discussed in Section 4.3 are enforced.

The one-dimensional functions $r_{n, j}$ are stored for all $n \in \mathbb{N}$ and $1 \leq j \leq d$ on a discrete grid with $M=30$ points, so that the minimization of functions of the form $\bigotimes_{j=1}^{d} r_{j} \mapsto$ $\mathcal{J}_{\nu_{t}}\left(f+\bigotimes_{j=1}^{d} r_{j}\right)$ is restricted to tensor products of one-dimensional continuous piecewise linear functions on this grid, and the Euler-Lagrange equations (17) are replaced by $M \times M$ linear systems. This discretization replaces the regularization by a kernel $K$, which is no more necessary. The process (3) is discretized with a time-step $\delta t=25.10^{-5}$, while the time integral in the empirical measure $\nu_{t}$ defined in (6) is discretized with a time-step $\Delta t=20 \delta t$. In other words, the reaction coordinates and the associated local mean forces are recorded in memory only every 20 steps of the Euler scheme. Moreover, $N$ independent replicas of the processes are run in parallel and the occupation measure used to defined the bias is $\tilde{\nu}_{N, t}$ given by (35). The update times $t_{k}$ of the bias are fixed at $t_{k}=k T$, with $T$ a multiple of $\Delta t$ and the number of tensor terms $g_{n}$ added at each update time is fixed with value $m$.

### 5.1 A low dimensional example

We start to test the method on a toy model, with $N=30$ replicas, a bias update period of $T=100 \Delta t$, a regularization parameter $\lambda=10^{-5}$, and $m=8$ tensor products added at each update. The reaction coordinates are Euclidian coordinates, more precisely $\xi(x)=\left(x_{1}, x_{2}\right)$, so that we don't introduce any additional extended coordinate.

The dimensions are $D=3, d=2$, particles start at $(0,0,0)$ and

$$
\begin{aligned}
V\left(x_{1}, x_{2}, x_{3}\right)= & -\sin \left(3 x_{1}\right) \sin \left(x_{2}\right) \cos \left(x_{3}-1\right)+\cos \left(3 x_{2}+2\right)\left(0.5+\cos \left(x_{3}-2\right)\right) \\
& +2 \sin \left(2 x_{1}+0.5\right) \cos \left(x_{3}\right)-5 \cos \left(x_{1}\right) \cos \left(x_{2}\right) \cos \left(x_{3}+1\right)
\end{aligned}
$$

This potential has the following properties: it is not a tensor product and yields a metastable process but, since $V\left(x_{1}, x_{2}, x_{3}\right)=\psi\left(x_{1}, x_{2}\right) \cos \left(x_{3}+\varphi\left(x_{1}, x_{2}\right)\right)$ for some functions $\psi$ and $\varphi$, there is no metastability in the orthogonal space for fixed $x_{1}, x_{2}$.

The results are given in Figures 1 and 2 (for $\beta=1$ ) and 3 and 4 (for $\beta=5$ ). In both cases, the theoretical free energy is successfully computed and the histograms of the reaction coordinates is eventually flat. This is a bit slower with the inverse temperature $\beta=5$, since the initial metastability is very strong. As can be seen in Figure 4, at that temperature and in the same times, a non-biased process is stuck in its initial well.

### 5.2 Polymer ring in solvent

We now consider a system inspired from [1]. The system is constituted of two types of particles, solvent particles, and polymer particles. The polymer particles interact through a potential made precise below to form a ring. The reaction coordinates are the bond lengths between consecutive polymer particles. This gives a large dimensional problem, for which the total dimension and the number of reaction coordinates are easily prescribed, and moreover where the reaction coordinates should exhibit some correlations (if it wasn't the case, then the TABF algorithm would not give better results than GABF [27]).

In a two-dimensional periodic box, we consider $D / 2=100$ particles among which $d$ (labeled from 1 to $d$ ) form a polymer and the others are solvent particles. The length of the box is $L=\sqrt{D / 2}$, to ensure a concentration independent from $D$. Each pair of particles that involves


Figure 1: For $\beta=1$, left and middle: estimated free energy respectively at $t=1$ and $t=30$. Right: theoretical free energy.



Figure 2: For $\beta=1$, cumulated histograms of the reaction coordinates at $t=30$ respectively for the TABF algorithm (left) and a non-biased process (right).


Figure 3: For $\beta=5$, left and middle: estimated free energy respectively at $t=10$ and $t=100$. Right: theoretical free energy.
at least one solvent particle interacts through the purely repulsive WCA pair potential, which



Figure 4: For $\beta=5$, cumulated histograms of the reaction coordinates at $t=100$ respectively for the TABF algorithm (left) and a non-biased process (right).
is the Lennard-Jones potential truncated at its minimum, namely

$$
V_{W C A}(r)=\varepsilon \mathbb{1}_{r \leqslant r_{0}}\left(1+\left(\frac{\sigma}{r}\right)^{12}-\left(\frac{\sigma}{r}\right)^{6}\right)
$$

where $r$ denotes the distance between the two particles, $\varepsilon=1, \sigma=0.5$ and $r_{0}=2^{1 / 6} \sigma$. Each pair of consecutive particles in the polymer ring (where the $d^{t h}$ and first particles are considered to be consecutive, closing the loop) interacts through a double well potential

$$
V_{D W}(r)=h\left(1-\frac{\left(2 r-2 r_{1}-\omega\right)^{2}}{\omega^{2}}\right)^{2}
$$

where $r_{1}=r_{0}, \omega=1$ and $h=3$. The minimum of $V_{D W}$ is attained at $r=r_{1}$ (compact state) and $r=r_{1}+w$ (stretched state). Finally, each triplet of consecutive particles in the polymer also interacts through the angle $\theta$ they form with the potential

$$
V_{A}(\theta)=\frac{1}{2}\left(\cos (\theta)-\cos \left(\theta_{d}\right)\right)^{2}
$$

with an equilibrium angle $\theta_{d}=\pi(1-2 / d)$ that ensures that the total angular potential is minimized when the polymer particles form a regular polygon.

There are $d$ reaction coordinates, which are the distances between two consecutive polymer particles. Following Section 4.1, the interaction between an extended reaction coordinate $z$ and the corresponding distance $r$ in the system is given via the extended potential

$$
V_{E}(z, r)=\frac{1}{2 \delta}\left(z-\frac{r-r_{1}}{w}\right)^{2}
$$

for $\delta=0.01$. The scaling ensures that the minimum of $V_{E}(z, r)+V_{D W}(r)$ is attained at $z=\left(r-r_{1}\right) / w \in\{0,1\}$. Moreover, in line with Section 4.2, the extended variable is confined in $\left[\xi_{\text {min }}, \xi_{\text {max }}\right]^{d}$ by orthogonal reflection at the boundary, with $\xi_{\text {min }}=-0.2$ and $\xi_{\max }=1.2$.

The total energy of the (extended) system is thus, for $(q, z) \in(L \mathbb{T})^{D} \times\left[\xi_{\text {min }}, \xi_{\text {max }}\right]^{d}$,

$$
\begin{aligned}
V(q, z)= & \sum_{i=d+1}^{D / 2} \sum_{j<i} V_{W C A}\left(\left|q_{i}-q_{j}\right|\right)+\sum_{i=1}^{d} V_{E}\left(z_{i},\left|\tilde{q}_{i+1}-q_{i}\right|\right) \\
& +\sum_{i=1}^{d} V_{D W}\left(\left|q_{i}-\tilde{q}_{i+1}\right|\right)+\sum_{i=1}^{d-1} V_{A}\left(\arccos \left(\frac{\tilde{q}_{i+1}-q_{i}}{\left|\tilde{q}_{i+1}-q_{i}\right|} \cdot \frac{\tilde{q}_{i+2}-\tilde{q}_{i-1}}{\left|\tilde{q}_{i+2}-\tilde{q}_{i+1}\right|}\right)\right),
\end{aligned}
$$



Figure 5: The slow motions of the system are the transitions of each bond between two consecutive particles of the polymer from its compact state to its stretched state.
where $\tilde{q}_{i}=q_{i}$ for all $i \in \llbracket 1, d \rrbracket$ and $\tilde{q}_{d+j}=q_{j}$ for $j=1,2$. Initially, the polymer is in a compact state, i.e. the distances between two consecutive of its particles are at distance $r_{1}$, the angles are $\theta_{d}$ and all the extended variables $\left(z_{i}\right)_{i \in \llbracket 1, d \rrbracket}$ are at 0 . For this model, we use the variant described in point 5 of Section 4.4 namely, following the GABF algorithm, we keep in memory one dimensional free energies on a grid and we use the tensor approximation as a correction of this initial guess. There are $N=50$ replicas, the update period, regularization parameter, and inverse temperature are respectively $T=10^{4} \Delta t, \lambda=0.05$ and $\beta=1$, and at each update, $m=4 d$ tensor products are added.

The free energy is expected to be close to a sum of one-dimensional double well potentials, with minima attained at points close to 0 and 1 . Nevertheless the angular force should favor configurations where the consecutive distances in the polymer are close. This fact cannot be grasped by the GABF algorithm alone, for which reaction coordinates are treated independently one from the others.

The results are presented in Figure 6 for $d=3$ and Figures 7 and 8 for $d=5$. In Figure 6 , we see that indeed the one-dimensional free energies recovered by the GABF algorithm have two wells approximately at 0 and 1 , and that the non-independent part of the free energy has the following effect: when $z_{3}=0$, the well $(0,0)$ is favored, when $z_{3}=1$ the same goes for $(1,1)$, and when $z_{3}$ is intermediate the landscape is flatter and the two wells $\left\{z_{1}=z_{2}=x\right\}$ with $x \in\{0,1\}$ are favored with respect to the wells $(0,1)$ and $(1,0)$. The result is similar in Figure 7, even though the quality of the estimation is lower for $z_{3}=z_{4}=z_{5}=0.5$, which is to be expected as this lies in a very low probability area (since 0.5 is the saddle point of the two well potential). This shows that the TABF algorithm is able to recover non-trivial correlations between reaction coordinates.

## References

[1] H. Alrachid and T. Lelièvre. Long-time convergence of an adaptive biasing force method: Variance reduction by helmholtz projection. SMAI Journal of Computational Mathematics, 1:55-82, 2015.
[2] L. Ambrosio, A. Carlotto, and A. Massaccesi. Lectures on elliptic partial differential equations, volume 18 of Appunti. Scuola Normale Superiore di Pisa (Nuova Serie). Edizioni della Normale, Pisa, 2018.


Figure 6: For $d=3$, up: estimated free energy as a function of $\left(z_{1}, z_{2}\right)$ when $z_{3}$ is, respectively, 0 (left) 0.5 (middle) and 1 (right). Middle: idem but without the independent, one-dimensional parts given by the GABF algorithm. Bottom: one-dimensional potential given by the GABF algorithm for $z_{1}$ (left) $z_{2}$ (middle) and their sum (right).
[3] D. Bakry, I. Gentil, and M. Ledoux. Analysis and geometry of Markov diffusion operators, volume 348 of Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. Springer, Cham, 2014.
[4] M. Benaïm and C.-E. Bréhier. Convergence of adaptive biasing potential methods for diffusions. C. R. Math. Acad. Sci. Paris, 354(8):842-846, 2016.
[5] M. Benaïm, C.-E. Bréhier, and P. Monmarché. Analysis of an Adaptive Biasing Force method based on self-interacting dynamics. arXiv e-prints, page arXiv:1910.04428, Oct 2019.
[6] M. Benaïm, M. Ledoux, and O. Raimond. Self-interacting diffusions. Probab. Theory Related Fields, 122(1):1-41, 2002.
[7] A. Bittracher, R. Banisch, and C. Schütte. Data-driven computation of molecular reaction coordinates. The Journal of Chemical Physics, 149(15):154103, 2018.
[8] S. Brandt, F. Sittel, M. Ernst, and G. Stock. Machine learning of biomolecular reaction coordinates. The Journal of Physical Chemistry Letters, 9(9):2144-2150, 2018.
[9] E. Cancès, V. Ehrlacher, and T. Lelièvre. Convergence of a greedy algorithm for highdimensional convex nonlinear problems. Math. Models Methods Appl. Sci., 21(12):24332467, 2011.


Figure 7: For $d=5$, up: estimated free energy as a function of $\left(z_{1}, z_{2}\right)$ when $z_{3}, z_{4}, z_{5}$ are, respectively, $(0,0,0)$ (left) $(0.5,0.5,0.5)$ (middle) and ( $1,1,1$ ) (right). Bottom: idem but without the independent, one-dimensional parts given by the GABF algorithm.


Figure 8: For $d=5$, cumulative one-dimensional histograms of the five reaction coordinates at $t=50$ for the TABF algorithm (up) and for a non biased process (bottom).
[10] E. Darve and A. Pohorille. Calculating free energies using average force. The Journal of Chemical Physics, 2001.
[11] L. De Lathauwer, B. De Moor, and J. Vandewalle. A multilinear singular value decomposition. SIAM J. Matrix Anal. Appl., 21(4):1253-1278, 2000.
[12] M. Espig, W. Hackbusch, and A. Khachatryan. On the Convergence of Alternating Least Squares Optimisation in Tensor Format Representations. arXiv e-prints, page arXiv:1506.00062, May 2015.
[13] H. Fu, X. Shao, C. Chipot, and W. Cai. Extended adaptive biasing force algorithm. an on-the-fly implementation for accurate free-energy calculations. J. Chem. Theory Comput., 12(8):3506-3513, 2016.
[14] L. Grasedyck, D. Kressner, and C. Tobler. A literature survey of low-rank tensor approximation techniques. GAMM-Mitteilungen, 36(1):53-78, 2013.
[15] J. Hénin and C. Chipot. Overcoming free energy barriers using unconstrained molecular dynamics simulations. The Journal of Chemical Physics, 121:2904-2914, 2004.
[16] S. Klus and C. Schütte. Towards tensor-based methods for the numerical approximation of the Perron-Frobenius and Koopman operator. Journal of Computational Dynamics, 3:139, 2016.
[17] C. Le Bris, T. Lelièvre, and Y. Maday. Results and questions on a nonlinear approximation approach for solving high-dimensional partial differential equations. Constr. Approx., 30(3):621-651, 2009.
[18] T. Lelièvre. Two mathematical tools to analyze metastable stochastic processes. In Numerical mathematics and advanced applications 2011, pages 791-810. Springer, Heidelberg, 2013.
[19] T. Lelièvre, M. Rousset, and G. Stoltz. Long-time convergence of an adaptive biasing force method. Nonlinearity, 21(6):1155-1181, 2008.
[20] T. Lelièvre, M. Rousset, and G. Stoltz. Free energy computations: A mathematical perspective. Imperial College Press, 2010.
[21] A. Lesage, T. Lelièvre, G. Stoltz, and J. Hénin. Smoothed biasing forces yield unbiased free energies with the extended-system adaptive biasing force method. The Journal of Physical Chemistry B, 121(15):3676-3685, 2017. PMID: 27959559.
[22] C. Lu, X. Li, D. Wu, L. Zheng, and W. Yang. Predictive sampling of rare conformational events in aqueous solution: Designing a generalized orthogonal space tempering method. Journal of Chemical Theory and Computation, 12(1):41-52, 2016.
[23] P. Monmarché. Generalized $\Gamma$ calculus and application to interacting particles on a graph. Potential Anal., 50(3):439-466, 2019.
[24] F. Nüske, R. Schneider, F. Vitalini, and F. Noé. Variational tensor approach for approximating the rare-event kinetics of macromolecular systems. The Journal of Chemical Physics, 144(5):054105, 2016.
[25] G. Pagès. Introduction to vector quantization and its applications for numerics. In CEMRACS 2013-modelling and simulation of complex systems: stochastic and deterministic approaches, volume 48 of ESAIM Proc. Surveys, pages 29-79. EDP Sci., Les Ulis, 2015.
[26] S. Piana and A. Laio. A bias-exchange approach to protein folding. The Journal of Physical Chemistry B, 111(17):4553-4559, 2007. PMID: 17419610.
[27] T. Zhao, H. Fu, T. Lelièvre, X. Shao, C. Chipot, and W. Cai. The extended generalized adaptive biasing force algorithm for multidimensional free-energy calculations. Journal of Chemical Theory and Computation, 13(4):1566-1576, 2017.

