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Stochastic Optimization with Unadjusted Kernel: the SOUK Algorithm

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Abstract

In this paper we propose a stochastic approximation algorithm to minimize functions for which the gradient writes as the expectation of some integrable features. Considering the case where computing feature moments or sampling from the underlying probability distribution is not feasible we mix a Markov chain dynamic with the stochastic gradient descent dynamic, thus combining optimization tools with statistics controls. This approach is motivated with examples from Empirical Bayesian statistics. We assess the convergence as well as rates of convergence for various objective functions.

1 Introduction

Computational Bayesian statistics often rely on complex probabilistic modeling which involves latent variables and therefore the underlying a posteriori distribution of the parameters is expressed as an untractable parameter-dependent integral. For example, some characteristics of the prior distribution such as the mean or the variance are frequently unknown and are treated as hyperparameters which leads to consider hierarchical Bayesian models [42]. A first approach is to consider the joint a posteriori probability distribution and perform inference for the parameters of interest by marginalizing along the hyperparameters. This framework is referred to as a fully Bayesian setting. However, solving this problem can be computationally prohibitive since it requires to sample from the joint probability distribution. For instance this can be realized with a Gibbs sampling strategy but the convergence of the associated Markov chain can be very slow in some situations.

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A second approach, which we consider in this paper, referred to as the Empirical Bayes (EB) setting [7, 8, 40], consists in first finding a single-point estimate for the hyperparameters, similarly to the Maximum A Posteriori (MAP) estimator, based on the conditional probability distribution of the hyperparameters given the observations. Second, this estimate is plugged in the a posteriori distribution of the parameters given the observations and the hyperparameters in order to make inference. This methodology is theoretically founded: [37, 28, 45, 11] show that fully Bayesian and EB approaches lead to similar inference conclusions when the sample size is large enough under additional appropriate conditions. However, EB inference remains a challenging problem since it relies on the optimization of a high-dimensional parameter-depending integral. When the parameter space is small, numerical integration techniques can be considered, like the recently introduced nested Laplace approximations [35, 46], however they might introduce some bias which can be hard to quantify. If the dimension is large, the optimization can be achieved using a Robbins-Monro scheme [41, 29, 20] which requires to be able to sample from the probability distribution of the parameters given the observations and the hyperparameters. This former step is usually tackled through Monte Carlo or Monte Carlo Markov Chain (MCMC) algorithms. EB inference has found many applications in Bayesian inference problem [21, 6], large-scale inverse problem [50], biostatistics [47, 26, 17].

In contrast to previous work, in this contribution we propose and analyze the use of inexact MCMC algorithms to get approximate samples from the a posteriori distribution of the parameters given the observations and the hyperparameters. While this method introduces some bias in the estimation of the stochastic gradient used in the Robbins-Monro procedure, we show that it can be controlled explicitly and does not affect the convergence of the hyperparameter recursion. The use of inexact MCMC algorithms is justified because while existing results can be applied to stochastic approximation based on exact MCMC schemes it requires strong conditions on the associated Markov kernel (see e.g. [2, H6]). The basic condition of geometric ergodicity can be hard to check, for example for the Metropolis Adjusted Langevin Algorithm [43] or Hamiltonian Monte Carlo [13]. On the other hand, inexact MCMC algorithms such that the Unadjusted Langevin-based Monte Carlo methods are easier to analyze and can exhibit stronger convergence properties [5, 12, 9, 16]. In addition, they give in practice similar inference results than their exact counterparts at the cost of a speed and accuracy trade-off. As a result, there have been increasingly popular in the computational statistics and machine learning communities [49, 33, 18, 22, 36, 24].

The paper is organized as follows. In Section 2 we motivate and introduce our new methodology based on a Robbins-Monro procedure using inexact MCMC algorithms for Empirical Bayes estimation. In Section 3 we summarize our main theoretical contributions and check that they hold for Langevin-based methods. In particular, we show that under appropriate conditions the complexity of the proposed methodology is linear in the dimension. Rigorous theoretical treatment of the convergence of the algorithm is given in the appendix.

Notations and convention

Denote by $\mathcal{B}(\mathbb{R}^d)$ the Borel σ -field of \mathbb{R}^d , $\mathbb{F}(\mathbb{R}^d)$ the set of all Borel measurable functions on \mathbb{R}^d and for $f \in \mathbb{F}(\mathbb{R}^d)$, $||f||_{\infty} = \sup_{x \in \mathbb{R}^d} |f(x)|$. For μ a probability measure on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$ and $f \in \mathbb{F}(\mathbb{R}^d)$ a μ -integrable function, denote by $\mu(f)$ the integral of f w.r.t. μ . For $f \in \mathbb{F}(\mathbb{R}^d)$, the V-norm of f is given by $||f||_V = \sup_{x \in \mathbb{R}^d} |f(x)|/V(x)$. Let ξ be a finite signed measure on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$. The V-total variation distance of ξ is defined as

$$\|\xi\|_V = \sup_{f \in \mathbb{F}(\mathbb{R}^d), \|f\|_V \leq 1} \left| \int_{\mathbb{R}^d} f(x) \mathrm{d}\xi(x) \right| \ .$$

If $V \equiv 1$, then $\|\cdot\|_V$ is the total variation denoted by $\|\cdot\|_{\text{TV}}$.

Let U an open set of \mathbb{R}^d . We denote by $C^k(U, \mathbb{R}^p)$, respectively $C^k_c(U, \mathbb{R}^p)$, the set of \mathbb{R}^p -valued k-differentiable functions, respectively the set of compactly supported \mathbb{R}^p -valued k-differentiable functions. Let $f \in C^1(U, \mathbb{R})$, we denote by $\nabla_x f$, the gradient of f. f is said to me m-convex with $m \ge 0$ if for all $x, y \in \mathbb{R}^d$

$$\langle \nabla_x f(x) - \nabla_x f(y), x - y \rangle \ge m \|x - y\|^2$$

We recall that if f is twice differentiable at point a, its Laplacian is given by $\Delta_x f(a) = \sum_{i=1}^d \frac{\partial^2 f}{\partial x_i^2}(a)$. Let $F \in C^1(U \times V, \mathbb{R})$ we denote by ∇_x , the gradient of F along the first coordinate and ∇_y the gradient of F along the second coordinate. For any $A \subset \mathbb{R}^d$, we denote by ∂A the boundary of A. Let $A \in \mathcal{B}(\mathbb{R}^d)$, we define $\operatorname{Vol}(A) = \operatorname{Leb}(\mathbb{1}_A)$ where Leb is the Lebesgue measure over $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$. Let $B \subset \mathbb{R}^d$, we introduce its convex hull $\operatorname{Conv}(B)$, the intersection of all the convex sets containing B. Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space. We denote by $L^2(\Omega, \mathcal{F}) = \{X : X \text{ is a random variable on } \Omega \text{ such } \mathbb{E}[X^2] < +\infty\}$. Denote by $\mu \ll \nu$ if μ is absolutely continuous w.r.t. ν and $d\mu/d\nu$ an associated density. Let μ, ν be two probability measures on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$.

$$\operatorname{KL}(\mu|\nu) = \begin{cases} \int_{\mathbb{R}^d} \frac{\mathrm{d}\mu}{\mathrm{d}\nu}(x) \log\left(\frac{\mathrm{d}\mu}{\mathrm{d}\nu}(x)\right) \mathrm{d}\nu(x) , & \text{if } \mu \ll \nu \\ +\infty & \text{otherwise} \end{cases}$$

2 EB inference based on inexact MCMC methods, the Stochastic Optimization with Unadjusted Kernel (SOUK)

2.1 EB inference

Consider the hierarchical model based on the observations y in the subspace $\mathbf{Y} \subset \mathbb{R}^{d_y}$ specified by the conditional density with respect to the Lebesgue measure $(y, x, \theta) \mapsto p(y|x, \theta)$, and prior distributions $(x, \theta) \mapsto p(x|\theta)$ and $\theta \mapsto p(\theta)$ where $x \in \mathbb{R}^d$ is the parameter of interest and $\theta \in \Theta \subset \mathbb{R}^m$ is a hyperparameter. Applying Bayes formula, we get that the a posteriori distribution is given for any $x \in \mathbb{R}^d$ by $p(x|y) \propto \int_{\Theta} p(y|x, \theta) p(x|\theta) p(\theta) d\theta$. The problem of sampling from this distribution is the subject of many works and remains a challenge. A common strategy relies on the Gibbs algorithm which consists in sampling from $p(\theta|x, y)$ and $p(x|y, \theta)$ alternatively and performing inference using the marginal distribution along the variable x. However, the two underlying problems are either computationally very expensive or even impossible to solve. In this paper, we consider the EB setting where the a posteriori distribution of x given y is approximated by $p(y|x, \theta^*)p(x|\theta^*)p(\theta^*)$ up to a normalizing constant with

$$\theta^* \in \arg\max p(\theta|y) \text{ with } p(\theta|y) \propto \int p(y|\theta, x) p(x|\theta) p(\theta) dx = \int p(x|y, \theta) dx .$$
 (1)

In order to realize this task, the main difficulty is to estimate θ^* . To do so, we use a stochastic approximation approach defining the recursion $(\theta_n)_{n\in\mathbb{N}}$ starting from $\theta_0 \in \Theta$ for any $n \in \mathbb{N}$ by

$$\theta_{n+1} = \Pi_{\Theta}(\theta_n - \delta_{n+1}\Delta_{\theta_n}) , \qquad (2)$$

where Π_{Θ} is the projection onto Θ and, for any $\theta \in \Theta$, Δ_{θ} is an estimator of $\nabla_{\theta} \log p(\theta|y)$, where ∇_{θ} is the gradient with respect to the θ variable. This estimator can be written in the following form

$$\nabla_{\theta} \log p(\theta|y) = \int \frac{\nabla_{\theta} p(x, y, \theta)}{p(x, y, \theta)} \pi_{\theta}(x) \mathrm{d}x , \qquad (3)$$

where for any $\theta \in \Theta$, π_{θ} stands for the conditional distribution $x \mapsto p(x|y,\theta)$. This methodology was already studied and proposed in numerous papers [2, 1] in the case where Δ_{θ} is unbiased, directly sampling from π_{θ} , or based on a Markov Chain Monte Carlo (MCMC) algorithm targeting π_{θ} . In this framework, for any $n \in \mathbb{N}$

$$\Delta_{\theta_n} = \frac{\nabla_{\theta} p(\theta_n)}{p(\theta_n)} + m_n^{-1} \sum_{k=1}^{m_n} \frac{\nabla_{\theta} p(X_n^k, y | \theta_n)}{p(X_k^n, y | \theta_n)} , \qquad (4)$$

where $(m_n)_{n \in \mathbb{N}}$ is a sequence of positive integers and for any $n \in \mathbb{N}$, $(X_k^n)_{k \in \{0,...,m_n\}}$ is either a sequence either i.i.d. random variables with distribution π_{θ_n} or a Markov Chain targeting this distribution.

In this contribution, we propose to use inexact MCMC algorithms which do not target exactly π_{θ} but a close distribution in a sense which will precised below. Indeed, there have been recently several major results on this family of methodologies showing that while they are biased, they benefit from explicit convergence properties (and subsampling strategies can be applied in big-data setting) contrary to most Metropolis-Hastings type algorithms [16, 12, 9]. These approximate MCMC schemes are based on stochastic continuous dynamics $(Y_t^{\theta})_{t\geq 0}$ for which the target distribution π_{θ} is invariant. Two fundamental examples are the Langevin dynamics solution of the following Stochastic Differential Equation (SDE)

$$\mathrm{d}Y_t^\theta = \nabla_x U_\theta(Y_t^\theta) \mathrm{d}t + \sqrt{2} \mathrm{d}B_t \;, \tag{5}$$

or the kinetic Langevin dynamics solution of

$$\mathrm{d} Y^{\theta}_t = V^{\theta}_t \;, \qquad \mathrm{d} V^{\theta}_t = \nabla_x U_{\theta}(Y^{\theta}_t) \mathrm{d} t - V^{\theta}_t \mathrm{d} t + \sqrt{2} \mathrm{d} B_t$$

where for any $\theta \in \Theta$, $U_{\theta} : x \mapsto \log \pi_{\theta}(x)$ and $(B_t)_{t \geq 0}$ is a standard *d*-dimensional Brownian motion. Under mild assumptions on U_{θ} , these two SDEs admit strong solutions for which π_{θ} and $\tilde{\pi}_{\theta} : (x, v) \mapsto \pi_{\theta}(x) \exp(-\|v\|^2/2)/(2\pi)^{d/2}$ are invariant probability measures. In addition, explicit convergence of $(Y_t^{\theta})_{t \geq 0}$ (respectively $(Y_t^{\theta}, V_t^{\theta})_{t \geq 0}$) to π_{θ} (respectively to $\tilde{\pi}_{\theta}$), in different metrics have been obtained [44, 15, 14]. However, sampling path solutions of these dynamics is in general not feasible. Therefore discretizations have to be used instead. In this paper we mainly focus on the Euler-Maruyama discretizations [44] of (5)

$$X_{k+1} = X_k - \gamma \nabla_x U_\theta(X_k) + \sqrt{2\gamma} Z_{k+1} .$$
(6)

where $\gamma > 0$ is a stepsize and $(Z_k)_{k \in \mathbb{N}^*}$ is a sequence of i.i.d *d*-dimensional zero-mean Gaussian random variables with covariance matrix identity. We have all the necessary tools to present our new algorithm. Our new methodology can be also applied to a broad class of stochastic optimization problems. Indeed, as mentioned previously, the EB estimation (1) can be casted in this framework, see (2). Therefore, we present it in full generality in the next section.

2.2 Stochastic Optimization with Unadjusted MCMC

The recursion (6) defines a Markov kernel $R_{\gamma,\theta}$. Under appropriate conditions, for any $\gamma \in (0, \bar{\gamma})$, and $\theta \in \Theta$ with $\bar{\gamma} > 0$, we can show that $R_{\gamma,\theta}$ admits an invariant distribution, see Lemma 16 in the appendix. This observation motivates the following stochastic approximation scheme.

Let Θ be a convex closed set in \mathbb{R}^m . Consider an objective function $f : \Theta \to \mathbb{R}$ which we want to minimize. We assume that f is continuously differentiable and that its gradient is given for any $\theta \in \Theta$ by

$$\nabla f(\theta) = \int_{\mathbb{R}^d} H_{\theta}(x) \pi_{\theta}(\mathrm{d}x) , \qquad (7)$$

where $(\theta, x) \mapsto H_{\theta}(x) \in C(\Theta \times \mathbb{R}^{d}, \mathbb{R})$, for any $\theta \in \Theta$, H_{θ} is π_{θ} -integrable and $(\pi_{\theta})_{\theta \in \Theta}$ is a family of probability distributions over $(\mathbb{R}^{d}, \mathcal{B}(\mathbb{R}^{d}))$. Note that in the EB setting, we have from (3) that $H_{\theta}(x) = \nabla_{\theta} p(x, y, \theta) / p(x, y, \theta)$. As emphasized in the previous section, we propose a stochastic approximation scheme which relies on biased estimates of $\nabla f(\theta)$ through a family of Markov kernels $\{K_{\gamma,\theta}, \gamma \in (0, \bar{\gamma}) \text{ and } \theta \in \Theta\}$, for $\bar{\gamma} > 0$, such that for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$, $K_{\gamma,\theta}$ admits an invariant probability distribution $\pi_{\gamma,\theta}$ on $(\mathbb{R}^{d}, \mathcal{B}(\mathbb{R}^{d}))$. We assume in addition that the bias associated to the use of this family of Markov kernels can be controlled w.r.t. to γ uniformly in θ , *i.e.* for example there exists C > 0 such that for all $\gamma \in (0, \bar{\gamma})$ with $\bar{\gamma} > 0$ and $\theta \in \Theta$, $\|\pi_{\gamma,\theta} - \pi_{\theta}\|_{\mathrm{TV}} \leq C\gamma^{\alpha}$ with $\alpha > 0$. In our applications $K_{\gamma,\theta}$ stands for $R_{\gamma,\theta}$ for any $\gamma \in (0, \bar{\gamma})$ with $\bar{\gamma} > 0$ and $\theta \in \Theta$, where $R_{\gamma,\theta}$ is associated with (6).

Let now $(\delta_n)_{n \in \mathbb{N}^*} \in (\mathbb{R}^*_+)^{\mathbb{N}^*}$, $(m_n)_{n \in \mathbb{N}} \in (\mathbb{N}^*)^{\mathbb{N}}$ be sequences of stepsizes and batch sizes which will be used to define the sequence relatively to the variable θ similarly to (2) and (4). Let $(\gamma_n)_{n \in \mathbb{N}} \in (\mathbb{R}^*_+)^{\mathbb{N}}$ be a sequence of stepsizes which will be used to get approximate samples from π_{θ_n} , similarly to (6). Starting from $X_0^0 \in \mathbb{R}^d$ and $\theta_0 \in \Theta$ we define on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, $(\{X_k^n : k \in \{0, \ldots, m_n\}\}, \theta_n)_{n \in \mathbb{N}}$ by the following recursion for $n \in \mathbb{N}$ and $k \in \{0, \ldots, m_n - 1\}$

 X_{k+1}^n is a Markov chain with kernel K_{γ_n,θ_n} and $X_0^n = X_{m_{n-1}}^{n-1}$ conditionally to \mathcal{F}_{n-1} if $n \ge 1$ (8)

$$\theta_{n+1} = \Pi_{\Theta} \left[\theta_n - \frac{\delta_{n+1}}{m_n} \sum_{k=1}^{m_n} H_{\theta_n}(X_k^n) \right] , \qquad (9)$$

where Π_{Θ} is the projection onto Θ and where \mathcal{F}_n is defined as follows for all $n \in \mathbb{N}^*$

$$\mathcal{F}_n = \sigma \left(\theta_0, \{ (X_k^{\ell})_{k \in \{0, \dots, m_\ell\}} : \ell \in \{0, \dots, n\} \} \right) , \tag{10}$$

where $\{(X_k^{\ell})_{k \in \{0,...,m_\ell\}} : \ell \in \{0,...,n\}\}$ is given by (8). Note that such a construction is always possible by Kolmogorov extension theorem [27, Theorem 5.16] by (9), for any $n \in \mathbb{N}$, θ_{n+1} is \mathcal{F}_{n-1} measurable. Then the sequence of approximate minimizers of f is given by $(\hat{\theta}_N)_{N \in \mathbb{N}^*}$ where for any $N \in \mathbb{N}^*$

$$\hat{\theta}_N = \left\{ \sum_{n=1}^N \delta_n \theta_n \right\} \middle/ \left\{ \sum_{n=1}^N \delta_n \right\} . \tag{11}$$

The pseudo-code associated with this method is given in Algorithm 1.

Under different sets of conditions on $f, H, (\delta_n)_{n \in \mathbb{N}}, (\gamma_n)_{n \in \mathbb{N}}$ and $(m_n)_{n \in \mathbb{N}}$ we obtain that $(\theta_N)_{N \in \mathbb{N}^*}$ converges a.s. to an element of $\arg \min_{\Theta} f$. In particular we consider the case where f is assumed to be convex, in Section 3. To summarize our results, we establish that if $(\gamma_n)_{n \in \mathbb{N}}$ and $(\delta_n)_{n \in \mathbb{N}}$ go to 0 sufficiently fast, $\mathbb{E}[f(\hat{\theta}_N)] - \min_{\Theta} f$ goes to 0 with a quantitative rate of convergence. In the case where $(\gamma_n)_{n \in \mathbb{N}}$ is held fixed, *i.e.* for all $n \in \mathbb{N}$, $\gamma_n = \gamma$, we show that while $\mathbb{E}[f(\hat{\theta}_N)]$ does not converge to 0, there exists $C, \alpha > 0$ such that $\limsup_{N \to +\infty} \mathbb{E}[f(\hat{\theta}_N)] - \min_{\Theta} f \leq C \gamma^{\alpha}$. In the case where f is non-convex we show that we can apply some results of stochastic approximation which imply that the sequence $(\theta_n)_{n \in \mathbb{N}}$ converges a.s. to a stationary point of the projected ordinary differential equation associated with ∇f and Θ . We postpone this result to the appendix, see Appendix A, since it involves a theoretical background which we think is out of the scope of the main document.

Algorithm 1 Stochastic optimization with unadjusted Markov kernel (SOUK)

1: function SOUK
2: Inputs:

$$\{K_{\gamma,\theta} : \gamma > 0, \theta \in \Theta\}, (\gamma_n)_{n \in \mathbb{N}}, (\delta_n)_{n \in \mathbb{N}}, (m_n)_{n \in \mathbb{N}}, N$$
3: Initialize:

$$X_0^0 \in \mathbb{R}^d \text{ and } \theta_0 \in \Theta$$
4: for $n \in \{0, \dots, N-1\}$ do
5: if $n \ge 1$ then
6: $X_0^n = X_{m_{n-1}}^{n-1}$
7: end if
8: for $k \in \{0, \dots, m_n - 1\}$ do
9: $X_{k+1}^n \sim K_{\gamma_n, \theta_n}(X_k^n, \cdot)$
10: end for
11: $\theta_{n+1} = \prod_{\Theta} \left[\theta_n - \frac{\delta_{n+1}}{m_n} \sum_{k=1}^{m_n} H_{\theta_n}(X_k^n)\right]$
12: end for
13: Outputs:
 $\hat{\theta}_N = \left\{\sum_{n=1}^N \delta_n \theta_n\right\} / \left\{\sum_{n=1}^N \delta_n\right\}$
14: end function

3 Stochastic Optimization with Unadjusted MCMC

3.1 Convex objective function

We consider the following assumptions on Θ and f, where $\Theta \subset \mathbb{R}^m$.

A1. Θ is a convex compact set and $\Theta \subset \overline{B}(0, M_{\Theta})$ with $M_{\Theta} > 0$.

A2. There exist an open set $U \subset \mathbb{R}^m$ and $L_f \ge 0$ such that $\Theta \subset U$ and $f \in C^1(U, \mathbb{R})$ is convex and for any $\theta_1, \theta_2 \in \Theta$

$$\left\|\nabla f(\theta_1) - \nabla f(\theta_2)\right\| \leq L_f \left\|\theta_1 - \theta_2\right\|.$$

Note that under A1 and A2, $\arg \min_{\Theta} f$ is non-empty. The following assumption ensures that the gradient of the objective function can be written as the expectation of some integrable function.

A3. For any $\theta \in \Theta$, there exist H_{θ} : $\mathbb{R}^d \to \mathbb{R}^m$ and a probability distribution π_{θ} on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$ satisfying that $(\theta, x) \mapsto H_{\theta}(x)$ is measurable, H_{θ} is π_{θ} -integrable and (7).

We impose some stability condition on the stochastic process $\{(X_k^n)_{k \in \{0,...,m_n\}} : n \in \mathbb{N}\}$ defined by (8) and that for any $\gamma \in (0, \bar{\gamma})$, with $\bar{\gamma} > 0$ and $\theta \in \Theta$ the iterates of $K_{\gamma,\theta}$ are close enough to π_{θ} after a sufficiently large number of iterations.

H1. (i) There exists a measurable function $V : \mathbb{R}^d \to [1, +\infty)$ and $A_1 \ge 0$ such that for any $n \in \mathbb{N}$

$$\mathbb{E}\left[V(X_{m_n}^n) \middle| X_0^0\right] \leqslant A_1 V(X_0^0) , \qquad \mathbb{E}\left[V(X_0^0)\right] < +\infty .$$
⁽¹²⁾

(ii) In addition, there exist $\bar{\gamma} > 0$, $A_2, A_3 \ge 0$, $B : \mathbb{R}^*_+ \to \mathbb{R}_+$ such that for any $\gamma \in (0, \bar{\gamma})$, $\theta \in \Theta$, $x \in \mathbb{R}^d$ and $n \in \mathbb{N}$

$$\left| \int_{\mathbb{R}^d} K^n_{\gamma,\theta}(x, \mathrm{d}y) H_{\theta}(y) - \pi_{\theta}(H_{\theta}) \right| \leqslant A_2 \rho^{n\gamma} V(x) + A_3 B(\gamma) .$$
(13)

This assumption does not impose that the Markov kernel $K_{\gamma,\theta}$, with $\gamma \in (0, \bar{\gamma})$ where $\bar{\gamma} > 0$ and $\theta \in \Theta$ is uniformly V-geometrically ergodic. However if an appropriate ergodicity condition holds for the family of kernels $\{K_{\gamma,\theta}, \gamma \in (0, \bar{\gamma}), \theta \in \Theta\}$, as in **H**5 below, then (13) can be written as $\|\pi_{\gamma,\theta}(H_{\theta}) - \pi_{\theta}(H_{\theta})\| \leq A_3 B(\gamma)$ for any $\gamma \in (0, \bar{\gamma})$ with $\bar{\gamma} > 0$ and $\theta \in \Theta$, where $\pi_{\gamma,\theta}$ is the unique invariant probability distribution of $K_{\gamma,\theta}$. We show that this condition holds in the case of the Langevin algorithm in Proposition 18.

Theorem 1. Assume A1, A2, A3. Let $\bar{\gamma} > 0$, $(\gamma_n)_{n \in \mathbb{N}}$, $(\delta_n)_{n \in \mathbb{N}^*}$ be sequences of non-increasing positive real numbers and $(m_n)_{n \in \mathbb{N}}$ be a sequence of positive integers satisfying $\sup_{n \in \mathbb{N}} \delta_n < 1/L_f$, $\sup_{n \in \mathbb{N}} \gamma_n < \bar{\gamma}$ and

$$\sum_{n=0}^{+\infty} \delta_{n+1} = +\infty , \qquad \sum_{n=0}^{+\infty} \delta_{n+1} B(\gamma_n) < +\infty , \qquad \sum_{n=0}^{+\infty} \delta_{n+1} / (m_n \gamma_n) < +\infty .$$
(14)

Let $\{(X_k^n)_{k \in \{0,...,m_n\}} : n \in \mathbb{N}\}$ be given by (8). Assume in addition that **H1** holds. Then the following statements hold:

- (a) $(\theta_n)_{n \in \mathbb{N}}$ defined by (9) converges a.s. to some $\theta^* \in \arg \min_{\Theta} f$;
- (b) furthermore, a.s. there exists $C \ge 0$ such that for any $n \in \mathbb{N}^*$

$$\left\{ \sum_{k=1}^{n} \delta_k f(\theta_k) \middle/ \sum_{k=1}^{n} \delta_k \right\} - \min_{\Theta} f \leqslant C \middle/ \left(\sum_{k=1}^{n} \delta_k \right) \ .$$

We start with the following technical lemma. Consider $(\eta_n)_{n \in \mathbb{N}}$ the sequence defined for any $n \in \mathbb{N}$ by

$$\eta_n = m_n^{-1} \sum_{k=1}^{m_n} \left\{ H_{\theta_n}(X_k^n) - \pi_{\theta_n}(H_{\theta_n}) \right\} , \qquad (15)$$

Lemma 2. Under the hypotheses of Theorem 1 we have for any $n \in \mathbb{N}^*$

$$\mathbb{E}\left[\left\|\eta_{n}\right\|\right] \leqslant (A_{1} \vee 1)A_{2} \frac{\rho^{-\bar{\gamma}} \mathbb{E}\left[V(X_{0}^{0})\right]}{\log(1/\rho)\gamma_{n}m_{n}} + A_{3}B(\gamma_{n}) .$$

Proof. Using the definition of $(\mathcal{F}_n)_{n \in \mathbb{N}}$, see (10), the Markov property, Lemma 11 and H1-(ii) we have for any $n \in \mathbb{N}$

$$\mathbb{E}\left[\|\eta_{n+1}\||\mathcal{F}_{n}\right] \leq (1/m_{n+1}) \sum_{k=1}^{m_{n+1}} |K_{\gamma_{n+1},\theta_{n+1}}^{k} H_{\theta_{n+1}}(X_{m_{n}}^{n}) - \pi_{\theta_{n+1}} \left(H_{\theta_{n+1}}\right)|$$

$$\leq m_{n+1}^{-1} \sum_{k=1}^{m_{n+1}} \left\{A_{2}\rho^{k\gamma_{n+1}}V(X_{m_{n}}^{n}) + A_{3}B(\gamma_{n+1})\right\}$$

$$\leq \frac{A_{2}\rho^{-\bar{\gamma}}V(X_{m_{n}}^{n})}{\log(1/\rho)\gamma_{n+1}m_{n+1}} + A_{3}B(\gamma_{n+1}),$$

where for the last inequality we use that $\sum_{k=1}^{+\infty} \rho^{\gamma k} \leq \rho^{-\bar{\gamma}}/[\log(1/\rho)\gamma]$. In a similar manner, we have

$$\mathbb{E}\left[\|\eta_0\| | X_0^0\right] \leqslant \frac{A_2 \rho^{-\gamma} V(X_0^0)}{\log(1/\rho) \gamma_0 m_0} + A_3 B(\gamma_0) \ .$$

We conclude using (12) in H1-(i).

We now turn to the proof of Theorem 1.

Proof. The proof is an application of [2, Theorem 2, Theorem 3]

(a) To apply [2, Theorem 2], it is enough to show that the following series converge a.s.

$$\sum_{n=0}^{+\infty} \delta_{n+1} \langle \Pi_{\Theta}(\theta_n - \delta_{n+1} \nabla f(\theta_n)), \eta_n \rangle , \quad \sum_{n=0}^{+\infty} \delta_{n+1} \|\eta_n\| , \quad \sum_{n=0}^{+\infty} \delta_{n+1}^2 \|\eta_n\|^2 .$$

where the sequence $(\eta_n)_{n\in\mathbb{N}}$ is defined for any $n\in\mathbb{N}$ by (15). Since $\Pi_{\Theta}(\theta_n - \delta_{n+1}\nabla f(\theta_n))$ is bounded, we are reduced to proving that a.s. $\sum_{n=0}^{+\infty} \delta_{n+1} ||\eta_n|| < +\infty$. As a result, using (14) and Lemma 2 we obtain that $\sum_{n\in\mathbb{N}} \delta_{n+1}\mathbb{E}[||\eta_n||] < +\infty$, which implies the stated convergence applying [2, Theorem 2].

(b) Combining [2, Theorem 3] with A1 we obtain that a.s. for any $n \in \mathbb{N}$

$$\sum_{k=1}^{n} \delta_{k} \left\{ f(\theta_{k}) - \min_{\Theta} f \right\}$$

$$\leq \frac{\|\theta_{0} - \theta^{*}\|^{2}}{2} - \sum_{k=0}^{n-1} \delta_{k+1} \langle \Pi_{\Theta}(\theta_{k} - \delta_{k+1} \nabla f(\theta_{k})) - \theta^{*}, \eta_{k} \rangle + \sum_{k=0}^{n-1} \delta_{k+1}^{2} \|\eta_{k}\|^{2}$$

$$\leq 2M_{\Theta}^{2} + 2M_{\Theta} \sum_{k=0}^{n-1} \delta_{k+1} \|\eta_{k}\| + \sum_{k=0}^{n-1} \delta_{k+1}^{2} \|\eta_{k}\|^{2} .$$
(16)
(16)
(16)
(17)

Since $\sum_{k \in \mathbb{N}} \delta_{k+1} \|\eta_k\|$ is finite a.s. so is $\sum_{k \in \mathbb{N}} \delta_{k+1}^2 \|\eta_k\|^2$. The proof is then completed upon dividing (16) by $\sum_{k=1}^n \delta_k$.

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In (8), $X_0^n = X_{m_{n-1}}^{n-1}$ for $n \in \mathbb{N}^*$. This procedure is referred to as warm-start in the sequel. In the proof of Theorem 1, X_0^n could be a random variable independent from \mathcal{F}_{n-1} for any $n \in \mathbb{N}$ with $\sup_{n \in \mathbb{N}^*} \mathbb{E}[V(X_0^n)] < +\infty$. This will not be the case in the fixed batch size case, see Theorem 4, where the warm-start procedure is crucial for the convergence to occur.

We extend this theorem to non convex objective function see Theorem 9 in Appendix A. Under the conditions of Theorem 1 with the additional assumption that $\partial \Theta$ is a smooth manifold we obtain that $(\theta_n)_{n \in \mathbb{N}}$ converges a.s. to some point θ^* such that $\nabla f(\theta^*) + \mathbf{n} = 0$ with $\mathbf{n} = 0$ if $\theta^* \in int(\Theta)$ and $\mathbf{n} \in T(\theta^*, \partial \Theta)^{\perp}$ if $\theta^* \in \partial \Theta$, where $T(\theta, \partial \Theta)$ is the tangent space of $\partial \Theta$ at point $\theta \in \partial \Theta$, see [3, Chapter 2].

In our applications $B(\gamma) = \gamma^{1/2}$, see Theorem 5. Assume then that there exist a, b, c > 0 such that for any $n \in \mathbb{N}^*$, $\delta_n = n^{-a}$, $\gamma_n = n^{-b}$ and $m_n = n^c$ then (14) is equivalent to

$$a < 1$$
, $a + b/2 > 1$, $a - b + c > 1$.

Suppose $a \in [0, 1)$ is given, then the previous equation reads

$$b = 2(1-a) + \varsigma_1$$
, $c = 3(1-a) + \varsigma_2$, $\varsigma_2 > \varsigma_1$. (18)

Thus b and c are proportional to 1-a. This illustrates a trade-off between the intrinsic inaccuracy of our algorithm through the family of Markov kernels (8) which do not exactly target π_{θ} and the minimization aim of our scheme. Note also that $(\delta_n)_{n \in \mathbb{N}}$ is allowed to be constant. This situation yield $\gamma_n = n^{-2-\varsigma_1}$ and $m_n = n^{3+\varsigma_2}$ with $\varsigma_2 > \varsigma_1 > 0$.

With the additional following assumption we establish quantitative bounds on $\mathbb{E}[f(\hat{\theta}_n) - \min_{\Theta} f]$, where $(\hat{\theta}_n)_{n \in \mathbb{N}}$ is given by (11).

H2. There exists $A_4 \ge 0$ such that for any $n \in \mathbb{N}$ we have

$$\mathbb{E}\left[\left(\sum_{k=1}^{m_n} H_{\theta_n}(X_k^n) - \pi_{\theta_n}(H_{\theta_n})\right)^2\right] \leqslant A_4 D(\gamma_n, m_n) \mathbb{E}\left[V(X_0^0)\right] .$$
(19)

Note that this assumption strengthens **H** 1 and imposes a control on the mean square error associated with the family of Markov kernels $\{K_{\gamma,\theta} : \gamma \in (0,\bar{\gamma}), \theta \in \Theta\}$ used to target $(\pi_{\theta})_{\theta \in \Theta}$ and the family of $\{x \mapsto H_{\theta}(x) : \theta \in \Theta\}$.

Theorem 3. Assume A1, A2, A3. Let $(\gamma_n)_{n \in \mathbb{N}}$, $(\delta_n)_{n \in \mathbb{N}}$ be sequences of non-increasing positive real numbers and $(m_n)_{n \in \mathbb{N}}$ be a sequence of positive integers satisfying $\sup_{n \in \mathbb{N}} \delta_n < 1/L_f$, $\sup_{n \in \mathbb{N}} \gamma_n < \bar{\gamma}$. Let $\{(X_k^n)_{k \in \{0,...,m_n\}} : n \in \mathbb{N}\}$ be given by (8). Assume in addition that H1 and H2 hold. Then, there exists $(E_n)_{n \in \mathbb{N}}$ such that

$$\mathbb{E}\left[\left\{\sum_{k=1}^{n} \delta_k f(\theta_k) \middle/ \sum_{k=1}^{n} \delta_k\right\} - \min_{\Theta} f\right] \leqslant E_n \middle/ \left(\sum_{k=1}^{n} \delta_k\right) ,$$

with for any $n \in \mathbb{N}$,

$$E_{n} = 2M_{\Theta}^{2} + 2M_{\Theta}\mathbb{E}\left[V(X_{0}^{0})\right] \left[\frac{(A_{1} \vee 1)A_{2}}{\rho^{\bar{\gamma}}\log(1/\rho)} \sum_{k=0}^{n} \{\delta_{k+1}/(\gamma_{k}m_{k})\} + A_{3}\sum_{k=0}^{n} \delta_{k+1}B(\gamma_{k})\right] + A_{4}\mathbb{E}\left[V(X_{0}^{0})\right] \sum_{k=0}^{n} \delta_{k+1}^{2}m_{k}^{-2}D(\gamma_{k},m_{k}).$$

$$(20)$$

Proof. Taking the expectation in (16) we get that for every $n \in \mathbb{N}$

$$\mathbb{E}\left[\sum_{k=0}^{n-1} \delta_{k+1}\left\{f(\theta_k) - \min_{\Theta} f\right\}\right] \leq 2M_{\Theta}^2 + 2M_{\Theta}\sum_{k=0}^{n-1} \delta_{k+1}\mathbb{E}\left[\|\eta_k\|\right] + \sum_{k=0}^{n-1} \delta_{k+1}^2\mathbb{E}\left[\|\eta_k\|^2\right] .$$
(21)

Using (15) and (19) we obtain that for any $n \in \mathbb{N}$

$$\mathbb{E}\left[\|\eta_n\|^2\right] \leqslant A_4 m_n^{-2} D(\gamma_n, m_n) \mathbb{E}\left[V(X_0^0)\right] .$$
(22)

Combining (21), Lemma 2 and (22) we conclude the proof.

In the case where $K_{\gamma,\theta} = R_{\gamma,\theta}$ is the Markov kernel associated with the Langevin update (6) for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$ with $\bar{\gamma} > 0$, $B(\gamma) = \gamma^{1/2}$ and $D(\gamma, m) = m^2$, see Theorem 5. Therefore if (14) and $\sum_{k=0}^{+\infty} \delta_{k+1}^2 < +\infty$ hold then $(E_n)_{n\in\mathbb{N}}$ is bounded by E_{∞} . Thus we obtain the L¹ convergence of the sequence $(f(\hat{\theta}_n))_{n\in\mathbb{N}}$ to $\min_{\Theta} f$ with rate $(\sum_{k=1}^n \delta_k)^{-1}$ and explicit bound E_{∞} . In the case of the Langevin update (6), if there exist a, b, c > 0 such that for any $n \in \mathbb{N}^*$, $\delta_n =$

In the case of the Langevin update (6), if there exist a, b, c > 0 such that for any $n \in \mathbb{N}^*$, $\delta_n = n^{-a}$, $\gamma_n = n^{-b}$ and $m_n = n^c$, (14) is equivalent to (18) and the accuracy, respectively the complexity, of the algorithm are of orders $(\sum_{k=1}^n \delta_k)^{-1} = \mathcal{O}(n^{a-1})$, respectively $\sum_{k=0}^n m_k = \mathcal{O}(n^{3(1-a)+\varsigma_2+1})$ for $\varsigma_2 > 0$. Thus for a fix target precision $\varepsilon > 0$ the complexity reads $\varepsilon^{-3} (\log(1/\varepsilon)/(1-a))^{1+\varsigma_2}$. On the other hand if we fix the complexity budget to N the accuracy is of order $N^{-(3+(1+\varsigma_2)/(1-a))^{-1}}$. These two considerations suggest to set a close to 0.

A case of interest is the fix stepsize setting, *i.e.* for all $n \in \mathbb{N}$, $\gamma_n = \gamma$. In that case if $(\delta_n)_{n \in \mathbb{N}}$ is non-increasing, $\sum_{n \in \mathbb{N}^*} \delta_n = +\infty$, $\lim_{n \to +\infty} m_n = +\infty$ and $\lim_{n \to +\infty} D(\gamma_0, m_n) m_n^{-2} \leq A_0(\gamma_0)$, using [38, Problem 80] we obtain that

$$\limsup_{n} \mathbb{E}\left[\left\{\sum_{k=1}^{n} \delta_{k} f(\theta_{k}) \middle/ \sum_{k=1}^{n} \delta_{k}\right\} - \min f\right] \leqslant \mathbb{E}\left[V(X_{0}^{0})\right] \left(M_{\Theta} A_{3} B(\gamma_{0}) + A_{4} \delta_{0} A_{0}(\gamma_{0})\right) \ .$$

Similar bounds can be obtained if $(m_n)_{n\in\mathbb{N}}$ and $(\gamma_n)_{n\in\mathbb{N}}$ are constant. However if $(m_n)_{n\in\mathbb{N}}$ is constant the convergence cannot be obtained using Theorem 1. Indeed in (14), the condition $\sum_{n\in\mathbb{N}} \delta_{n+1}\gamma_n^{-1} < +\infty$ is in contradiction with the condition $\sum_{n\in\mathbb{N}^*} \delta_n = +\infty$.

3.2 Convex objective function – Fixed batch size

In this section, we consider the case where the batch size is fixed, *i.e.* $m_n = m_0$ for all $n \in \mathbb{N}$. For ease of exposition we only consider the case where $m_0 = 1$. However the general case can be adapted from the proof of the result stated below. For simplicity of notation we let $\tilde{X}_{n+1} = X_1^n$ for any $n \in \mathbb{N}$. More precisely the recursion (8) and (9) can be written for any $n \in \mathbb{N}$ as

$$X_{n+1}$$
 is a sample from $K_{\gamma_n,\tilde{\theta}_n}(X_n,\cdot)$ conditionally to \mathcal{F}_n , (23)

$$\tilde{\theta}_{n+1} = \Pi_{\Theta} \left[\tilde{\theta}_n - \delta_{n+1} H_{\tilde{\theta}_n}(\tilde{X}_{n+1}) \right] , \qquad (24)$$

starting from $\theta_0 \in \mathbb{R}^d$, \tilde{X}_0 a random variable and where $\tilde{\mathcal{F}}_n$ is given by

$$\hat{\mathcal{F}}_n = \sigma\left(\hat{\theta}_0, (\hat{X}_\ell)_{\ell \in \{0, \dots, n\}}\right) .$$
⁽²⁵⁾

We consider the following additional assumptions. We impose some Lipschitz-regularity over $(\theta, x) \mapsto H_{\theta}(x)$.

A 4. There exists a measurable function $V : \mathbb{R}^d \to [1, +\infty)$ and $M_H, L_H \ge 0$ such that for any $x \in \mathbb{R}^d$ and $\theta_1, \theta_2 \in \Theta$ we have the following inequalities

$$||H_{\theta_1}(x) - H_{\theta_2}(x)|| \leq L_H ||\theta_1 - \theta_2||V(x)^{1/2}, \qquad ||H_{\theta_1}(x)|| \leq M_H V^{1/4}(x).$$

H3. There exists a measurable function $V : \mathbb{R}^d \to [1, +\infty)$ and $A_5 \ge 0$ such that for any $n \in \mathbb{N}$

$$\mathbb{E}\left[V(\tilde{X}_n)\big|\tilde{X}_0\right] \leqslant A_5 V(\tilde{X}_0) , \qquad \mathbb{E}\left[V(\tilde{X}_0)\right] < +\infty .$$

We consider a similar property as A4 on the family of kernels $\{K_{\gamma,\theta}, \gamma \in (0, \bar{\gamma}), \theta \in \Theta\}$, with $\bar{\gamma} > 0$, which weakens the assumption [2, H6].

H4. There exist $A_6 \ge 0$, $\Lambda_1 : (\mathbb{R}^*_+)^2 \to \mathbb{R}_+$ and $\Lambda_2 : (\mathbb{R}^*_+)^2 \to \mathbb{R}_+$ such that for any $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\bar{\gamma} > 0$, $\theta_1, \theta_2 \in \Theta$, $x \in \mathbb{R}^d$ and $a \in [1/4, 1/2]$

$$\|\delta_{x}K_{\gamma_{1},\theta_{1}} - \delta_{x}K_{\gamma_{2},\theta_{2}}\|_{V^{a}} \leqslant A_{6} \left[\Lambda_{1}(\gamma_{1},\gamma_{2}) + \Lambda_{2}(\gamma_{1},\gamma_{2})\|\theta_{1} - \theta_{2}\|\right] V^{2a}(x)$$

Finally, we suppose that each kernel $K_{\gamma,\theta}$ admits an invariant probability measure and converges geometrically to this probability distribution in V-norm with constants not depending on γ and θ . We show in [10] that this condition holds for the Langevin Monte Carlo algorithm.

H5. For all $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$ with $\bar{\gamma} > 0$, $K_{\gamma,\theta}$ is uniformly V-geometrically ergodic, i.e. there exists a probability distribution $\pi_{\gamma,\theta}$ on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$, $A_7 \ge 0$ and $\kappa \in (0, 1)$ such that for any $n \in \mathbb{N}$ and $a \in [1/4, 1]$

$$||K_{\gamma,\theta}^n - \pi_{\gamma,\theta}||_{V^a} \leqslant A_7 \kappa^{n\gamma} V^a(x) .$$

The following theorem ensures convergence properties for the algorithm similar to Theorem 1. The proof of this result is based on a generalization of [20, Lemma 4.2], which only holds in the case of exact MCMC schemes.

Theorem 4. Assume A1, A2, A3 and A4. Let $\bar{\gamma} > 0$, $(\gamma_n)_{n \in \mathbb{N}}$ and $(\delta_n)_{n \in \mathbb{N}^*}$ be sequences of non-increasing positive real numbers satisfying $\sup_{n \in \mathbb{N}} \delta_n < 1/L_f$, $\sup_{n \in \mathbb{N}} \gamma_n < \bar{\gamma}$, $\sup_{n \in \mathbb{N}} |\delta_{n+1} - \delta_n|\delta_n^{-2} < +\infty$, $\sum_{n=0}^{+\infty} \delta_{n+1} = +\infty$ and

$$\sum_{n=0}^{+\infty} \delta_{n+1} B(\gamma_n) < +\infty , \ \sum_{n=0}^{+\infty} \frac{\delta_{n+1}^2}{\gamma_n^2} < +\infty , \ \sum_{n=0}^{+\infty} \frac{\delta_{n+1}}{\gamma_{n+1}} \left[\Lambda_1(\gamma_n, \gamma_{n+1}) + \delta_{n+1} \Lambda_2(\gamma_n, \gamma_{n+1}) \right] < +\infty .$$
(26)

Let $(\tilde{X}_n)_{n \in \mathbb{N}}$ be given by (23). Assume in addition that $H_{1-(ii)}$, H_3 , H_4 , H_5 hold. Then the following statements hold:

- (a) $(\tilde{\theta}_n)_{n \in \mathbb{N}}$ defined by (24) converges a.s. to some $\theta^* \in \arg\min_{\Theta} f$;
- (b) furthermore, a.s. there exists $C \ge 0$ such that for any $n \in \mathbb{N}^*$

$$\left\{ \sum_{k=1}^{n} \delta_k f(\tilde{\theta}_k) \middle/ \sum_{k=1}^{n} \delta_k \right\} - \min_{\Theta} f \leqslant C \middle/ \left(\sum_{k=1}^{n} \delta_k \right) .$$

Proof. The proof of this theorem is postponed to Appendix B.1.

In our case of application $B(\gamma) = \gamma^{1/2}$, $\Lambda_1(\gamma_1, \gamma_2) = (\gamma_2^{-1} - \gamma_1^{-1})^{1/2} + \gamma_2^{-1/2} |\gamma_1 - \gamma_2|$ and $\Lambda_2(\gamma_1, \gamma_2) = \gamma_2^{1/2}$. Thus we obtain that the following series should be convergent

$$\sum_{n=0}^{+\infty} \delta_{n+1} \gamma_n^{1/2} < +\infty , \quad \sum_{n=0}^{+\infty} \delta_{n+1}^2 / \gamma_n^2 < +\infty , \quad \sum_{n=0}^{+\infty} \delta_{n+1} (\gamma_{n+1}^{-1} - \gamma_n^{-1})^{1/2} / \gamma_n^2 < +\infty .$$

A possible choice is given $\delta_n = n^{-1}$ and $\gamma_n = \log(n+1)^{-4}$. In this setting it should be noted that the sequence of $(\tilde{\theta}_n)_{n \in \mathbb{N}}$ moves slowly compared to the sequence $(\tilde{X}_n)_{n \in \mathbb{N}}$. In the limit case, if we suppose that $\tilde{\theta}_n$ is constant we obtain that averaging the sequence $(\tilde{X}_n)_{n \in \mathbb{N}}$ is an estimator of $\nabla f(\tilde{\theta}_n)$. This is coherent with the increasing batch size approach.

3.3 SOUL method

We consider the following assumption on the family of probability distributions $(\pi_{\theta})_{\theta \in \Theta}$.

L1. For any $\theta \in \Theta$, there exists $U_{\theta} : \mathbb{R}^{d} \to \mathbb{R}$ such that π_{θ} admits a probability density function w.r.t. to the Lebesgue measure proportional to $x \mapsto \exp(-U_{\theta}(x))$. In addition $(\theta, x) \mapsto U_{\theta}(x)$ is continuous, $x \mapsto U_{\theta}(x)$ is differentiable for all $\theta \in \Theta$ and there exists $L \ge 0$ such that for any $x, y \in \mathbb{R}^{d}$,

$$\sup_{\theta \in \Theta} \|\nabla_x U_\theta(x) - \nabla_x U_\theta(y)\| \leq \mathbf{L} \|x - y\| ,$$

and $\{\|\nabla_x U_{\theta}(0)\| : \theta \in \Theta\}$ is bounded.

Under L1, the Langevin diffusion defined by (5) admits a unique strong solution for any $\theta \in \Theta$. Therefore we consider the family of Markov kernels $\{R_{\gamma,\theta}, \gamma \in (0, \bar{\gamma}), \theta \in \Theta\}$ with $\bar{\gamma} > 0$, induced by the recursion (6). The recursion step in (8) can be rewritten for any $n \in \mathbb{N}$ and $k \in \{0, \ldots, m_n - 1\}$

$$X_{k+1}^{n} = X_{k}^{n} - \gamma_{n} \nabla_{x} U_{\theta_{n}}(X_{k}^{n}) + \sqrt{2\gamma_{n}} Z_{k+1}^{n} , \text{ with } X_{0}^{n} = X_{m_{n-1}}^{n-1} \text{ if } n \ge 1 , \qquad (27)$$

given a sequence of stepsizes $(\gamma_n)_{n \in \mathbb{N}}$, $(m_n)_{n \in \mathbb{N}} \in (\mathbb{N}^*)^{\mathbb{N}}$ and $(Z_k^n)_{n \in \mathbb{N}, k \in \{1, \dots, m_n\}}$ a family of i.i.d *d*-dimensional zero-mean Gaussian random variables with covariance matrix identity. Consider one of the following additional tail conditions on U_{θ} to ensure geometric ergodicity of $R_{\gamma,\theta}$ for any $\gamma \in (0, \bar{\gamma})$ and $\theta \in \Theta$, with $\bar{\gamma} > 0$ which will be specified below.

L2. There exist $\eta > 0$ and $m, c, M_{\eta} \ge 0$ such that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$,

$$\langle \nabla_x U_\theta(x), x \rangle \ge \eta \|x\| \mathbb{1}_{\overline{\mathrm{B}}(0,M_n)^{\mathrm{c}}}(x) + m \|\nabla_x U_\theta(x)\|^2 - c$$

L3. There exist $\eta > 0$ and $c, M_{\eta} \ge 0$ such that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$

$$\langle \nabla_x U_\theta(x), x \rangle \ge \eta \|x\|^2 \mathbb{1}_{\overline{\mathrm{B}}(0, M_n)^c}(x) - c ,$$

The next theorem asserts the convergence of the scheme defined by (9) and (27) referred to as the Stochastic Optimization with Unadjusted Langevin (SOUL) algorithm.

Theorem 5. Assume A1, A2, A3 and there exists $M_H \ge 0$ such that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $\|H_{\theta}(x)\| \le M_H V^{1/2}(x)$. In addition, assume L1 and L2 or L3. There exists $\bar{\gamma} > 0$ such that for all sequences of non-increasing positive real numbers $(\gamma_n)_{n\in\mathbb{N}}$, $(\delta_n)_{n\in\mathbb{N}^*}$ and sequence of positive integers $(m_n)_{n\in\mathbb{N}}$, satisfying $\sup_{n\in\mathbb{N}} \gamma_n < \bar{\gamma}$, H1 and H2 hold with $B(\gamma) = \gamma^{1/2}$ and $D(\gamma, m) = m^2$. *Proof.* The proof is postponed to Appendix B.2.4.

Consider the following assumption for $m \ge 0$.

L4 (m). For any $\theta \in \Theta$, U_{θ} is m-convex, continuously differentiable and satisfies

$$\sup_{\theta \in \Theta} \int_{\mathbb{R}^d} \exp(-U_{\theta}(x)) \mathrm{d}x < +\infty \; .$$

In addition, $\{y \in \mathbb{R}^d : \text{ there exists } \theta \in \Theta, \nabla_x U_\theta(y) = 0\} \subset B(0, M_*) \text{ with } M_* \ge 0.$

Note that $\{y \in \mathbb{R}^d : \text{there exists } \theta \in \Theta, \nabla_x U_\theta(y) = 0\}$ is the set of minimizers of the family of functions $(U_\theta)_{\theta \in \Theta}$.

Proposition 6. Assume A_1 , L_1 and $L_4(0)$. Then L_2 holds.

Proof. The proof is postponed to Appendix B.3.

While Proposition 6 shows that L4(0) and additional assumptions implies L2, it does not give an explicit expression for the constants appearing in this latter condition. In the case where L4(m)with m > 0 holds, we get the following result which implies that the bound provided by Theorem 3 is linear in the dimension.

Corollary 7. Assume A1, A2, A3 and there exists $M_H \ge 0$ such that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $\|H_{\theta}(x)\| \le M_H V^{1/2}(x)$. In addition assume L1 and L4(m), with m > 0. There exists $\bar{\gamma} > 0$ such that for all sequences of non-increasing positive real numbers $(\gamma_n)_{n\in\mathbb{N}}$, $(\delta_n)_{n\in\mathbb{N}^*}$ and sequence of positive integers $(m_n)_{n\in\mathbb{N}}$, satisfying $\sup_{n\in\mathbb{N}} \delta_n < 1/L_f$, $\sup_{n\in\mathbb{N}} \gamma_n < \bar{\gamma}$ and (14) with $B(\gamma) = \gamma^{1/2}$, and $D(\gamma, m) = m^2$, the conclusions of Theorem 3 hold with $(E_n)_{n\in\mathbb{N}}$ such that for any $n \in \mathbb{N}$, $|E_n| \le Cd$, where $C \ge 0$ is a constant independent of the dimension d.

Proof. The proof is postponed to Appendix B.4.

L5. There exists $L_U \ge 0$ such that for any $x \in \mathbb{R}^d$ and $\theta_1, \theta_2 \in \Theta$

$$\|\nabla_x U_{\theta_1}(x) - \nabla_x U_{\theta_2}(x)\| \leq L_U \|\theta_1 - \theta_2\|V(x)^{1/2}$$

Theorem 8. Assume A1, A2, A3 and A4. In addition assume L1, L5 and L2 or L3. There exists $\bar{\gamma} > 0$ such that for all sequences of non-increasing positive real numbers $(\gamma_n)_{n \in \mathbb{N}}$, $(\delta_n)_{n \in \mathbb{N}^*}$ and sequence of positive integers $(m_n)_{n \in \mathbb{N}}$, satisfying $\sup_{n \in \mathbb{N}} \gamma_n < \bar{\gamma}$, H4 and H5 hold with

$$\Lambda_1(\gamma_1,\gamma_2) = (\gamma_2^{-1} - \gamma_1^{-1})^{1/2} + \gamma_2^{-1/2} |\gamma_1 - \gamma_2| , \qquad \Lambda_2(\gamma_1,\gamma_2) \leqslant \gamma_2^{1/2} ,$$

Proof. The proof is postponed to Appendix B.5.

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A Non-convex objective function

We now turn to the case where f is non-convex. We recall that the normal space of a sub-manifold $\mathcal{M} \subset \mathbb{R}^m$ at point x is given by

$$N(x, \mathcal{M}) = \begin{cases} T(x, \mathcal{M})^{\perp} & \text{if } x \in \mathcal{M} ;\\ \{0\} & \text{otherwise }, \end{cases}$$

where $T(x, \mathcal{M})$ is the tangent space of the sub-manifold \mathcal{M} at point x, see [3].

Theorem 9. Assume A1 and that $\partial \Theta$ is a \mathbb{R}^{m-1} connected manifold with continuously differentiable outer normal. In addition, assume that $f \in C^1(\mathbb{R}^p, \mathbb{R})$, A3. Let $\bar{\gamma} > 0$, $(\gamma_n)_{n \in \mathbb{N}}$, $(\delta_n)_{n \in \mathbb{N}^*}$ be sequences of non-increasing positive real numbers and $(m_n)_{n \in \mathbb{N}}$ be a sequence of positive integers such that $\sup_{n \in \mathbb{N}} \delta_n < 1/L_f$, $\sup_{n \in \mathbb{N}} \gamma_n < \bar{\gamma}$ and (14) are satisfied. Let $\{(X_k^n)_{k \in \{0,...,m_n\}} : n \in \mathbb{N}\}$ be given by (8). Assume in addition H1 holds. Then $(\theta_n)_{n \in \mathbb{N}}$ defined by (9) converges a.s. to some $\theta^* \in \{\theta \in \Theta : \nabla f(\theta) + \mathbf{n} = 0, \mathbf{n} \in \mathbb{N}(\theta, \partial \Theta)\}$, where $\mathbb{N}(\theta, \partial \Theta)$ is the normal space of $\partial \Theta$ at point θ .

Theorem 10. Assume A1 and that $\partial \Theta$ is a \mathbb{R}^{m-1} connected manifold with continuously differentiable outer normal. In addition, assume that $f \in C^1(\mathbb{R}^p, \mathbb{R})$, A3 and A4. Let $\bar{\gamma} > 0$, $(\gamma_n)_{n\in\mathbb{N}}, (\delta_n)_{n\in\mathbb{N}^*}$ be sequences of non-increasing positive real numbers and $(m_n)_{n\in\mathbb{N}}$ be a sequence of positive integers such that $\sup_{n\in\mathbb{N}} \delta_n < 1/L_f$, $\sup_{n\in\mathbb{N}} \gamma_n < \bar{\gamma}$, $\sup_{n\in\mathbb{N}} |\delta_{n+1} - \delta_n| \delta_n^{-2} < +\infty$, $\sum_{n=0}^{+\infty} \delta_{n+1} = +\infty$ and (26) are satisfied. Let $(\tilde{X}_n)_{n\in\mathbb{N}}$ be given by (23). Assume in addition that H 1-(ii), H 3, H 4 and H 5 hold. Then $(\tilde{\theta}_n)_{n\in\mathbb{N}}$ defined by (24) converges a.s. to some $\theta^* \in \{\theta \in \Theta : \nabla f(\theta) + \mathbf{n} = 0, \mathbf{n} \in \mathbb{N}(\theta, \partial\Theta)\}$, where $\mathbb{N}(\theta, \partial\Theta)$ is the normal space of $\partial\Theta$ at point θ .

Proof. The proof is a direct application of [29, Chapter 5, Theorem 2.3] using the decomposition of the error term considered in the proof of Theorem 1 and Theorem 4. Indeed we decompose the error term η_n defined by (15) as $\eta_n = \delta M_n + B_n$, where δM_n is a martingale increment. Then, we only need to show that the following sums converge

$$\sum_{k=0}^{n} \delta_{k+1}^{2} \mathbb{E} \left[\|\delta M_{k}\|^{2} \right] , \qquad \sum_{k=0}^{n} \delta_{k+1} \mathbb{E} \left[\|B_{k}\| \right] .$$

This is exactly what we show in Theorem 1 with $\delta M_n = 0$ and $B_n = \eta_n$ and in Theorem 4 with $\delta M_n = \tilde{\eta}_n^a$ and $B_n = \eta_n^b + \eta_n^c + \eta_n^d$ where $\{\tilde{\eta}_n^i : n \in \mathbb{N}, i \in \{a, b, c, d\}\}$ is defined by (34).

B Postponed proofs

B.1 Proof of Theorem 4

In all the proof we use the following lemma

Lemma 11. Let $t \in (0,1)$ and $\gamma \in (0,\bar{\gamma})$ with $\bar{\gamma} > 0$ then $\sum_{n \in \mathbb{N}} t^{n\gamma} \leq t^{-\bar{\gamma}} \log^{-1}(1/t) \gamma^{-1}$ and $\sum_{n \in \mathbb{N}} nt^{n\gamma} \leq t^{-\bar{\gamma}} \log^{-2}(1/t) \gamma^{-2}$.

Proof. Let $t \in (0,1)$ and $\gamma \in (0,\bar{\gamma})$ with $\bar{\gamma} > 0$. Using that $e^u - 1 \leq ue^u$ for all $u \geq 0$ we have

$$\sum_{n \in \mathbb{N}} t^{n\gamma} = -(t^{\gamma} - 1)^{-1} \leqslant -\gamma^{-1} \log^{-1}(t) \exp(-\log(t)\gamma) \leqslant t^{-\bar{\gamma}} \log^{-1}(1/t)\gamma^{-1}$$

and

$$\sum_{n \in \mathbb{N}} n t^{n\gamma} = t^{\gamma} (t^{\gamma} - 1)^{-2} \leqslant t^{\gamma} \{ \gamma^{-1} \log^{-1}(t) \exp(-\log(t)\gamma) \}^2 \leqslant t^{-\bar{\gamma}} \log^{-2}(1/t)\gamma^{-2} ,$$

which completes the proof.

Consider $(\tilde{X}_k, \tilde{\theta}_k)_{k \in \mathbb{N}}$ defined by (23) and (24) and define for any $n \in \mathbb{N}$,

$$\tilde{\eta}_n = H_{\tilde{\theta}_n}(\tilde{X}_{n+1}) - \nabla f(\tilde{\theta}_n) = H_{\tilde{\theta}_n}(\tilde{X}_{n+1}) - \pi_{\tilde{\theta}_n}(H_{\tilde{\theta}_n})$$

In order to apply [2, Theorem 2] we show that the following sums converge a.s.

$$\sum_{k=0}^{n} \delta_{k+1} \langle \Pi_{\Theta} \left[\tilde{\theta}_{k} - \delta_{k} \nabla f(\tilde{\theta}_{k}) \right], \tilde{\eta}_{k} \rangle, \qquad \sum_{k=0}^{n} \delta_{k+1} \tilde{\eta}_{k}, \qquad \sum_{k=0}^{n} \delta_{k+1}^{2} \| \tilde{\eta}_{k} \|^{2}.$$
(28)

Using $||x+y||^2 \leq 2(||x||^2 + ||y||^2)$ for any $x, y \in \mathbb{R}^d$, A1, A2, A3, H3 and A4, we get for any $k \in \mathbb{N}$, $\mathbb{E}[||\tilde{\eta}_k||^2] \leq 2M_H^2 A_6 \mathbb{E}[V(\tilde{X}_0)] + 2\sup_{\theta \in \Theta} ||\nabla f(\theta)||^2 < +\infty$, which implies that for any $n \in \mathbb{N}$,

$$\mathbb{E}\left[\sum_{k=0}^{n} \delta_{k+1}^{2} \|\tilde{\eta}_{k}\|^{2}\right] \leqslant \sup_{k \in \mathbb{N}} \mathbb{E}\left[\|\tilde{\eta}_{k}\|^{2}\right] \sum_{k=1}^{n} \delta_{k}^{2} \leqslant \sup_{k \in \mathbb{N}} \mathbb{E}\left[\|\tilde{\eta}_{k}\|^{2}\right] \sum_{k=1}^{+\infty} \delta_{k}^{2} < +\infty .$$
(29)

By A4 and H5, for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$, there exists a function $\hat{H}_{\gamma,\theta} : \mathbb{R}^d \to \mathbb{R}^d$ solution of the *Poisson equation*,

$$(\mathrm{Id} - K_{\gamma,\theta})\dot{H}_{\gamma,\theta} = H_{\theta} - \pi_{\gamma,\theta}(H_{\theta})$$
(30)

,

defined for any $x \in \mathbb{R}^d$ by

$$\hat{H}_{\gamma,\theta}(x) = \sum_{j \in \mathbb{N}} \{ K^j_{\gamma,\theta} H_{\theta}(x) - \pi_{\gamma,\theta}(H_{\theta}) \} , \qquad (31)$$

which satisfies using for all $x \in \mathbb{R}^d$

$$\|\hat{H}_{\theta}(x)\| \leqslant C_P V^{1/4}(x) , \qquad (32)$$

where $C_P = A_7 \log(1/\kappa)^{-1} \kappa^{-\bar{\gamma}} \gamma^{-1}$ by Lemma 11. Using (30), we consider the following decomposition: for any $k \in \mathbb{N}$

$$\tilde{\eta}_{k} = \hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k+1}) - K_{\gamma_{k},\tilde{\theta}_{k}}\hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k+1}) + \pi_{\gamma_{k},\tilde{\theta}_{k}}(H_{\tilde{\theta}_{k}}) - \pi_{\tilde{\theta}_{k}}(H_{\tilde{\theta}_{k}}) = \tilde{\eta}_{k}^{a} + \tilde{\eta}_{k}^{b} + \tilde{\eta}_{k}^{c} + \tilde{\eta}_{k}^{d} ,$$
(33)

where

$$\begin{cases} \tilde{\eta}_{k}^{a} = \hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k+1}) - K_{\gamma_{k},\tilde{\theta}_{k}}\hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k}) ;\\ \tilde{\eta}_{k}^{b} = K_{\gamma_{k},\tilde{\theta}_{k}}\hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k}) - K_{\gamma_{k+1},\tilde{\theta}_{k+1}}\hat{H}_{\tilde{\theta}_{k+1}}(\tilde{X}_{k+1}) ;\\ \tilde{\eta}_{k}^{c} = K_{\gamma_{k+1},\tilde{\theta}_{k+1}}\hat{H}_{\tilde{\theta}_{k+1}}(\tilde{X}_{k+1}) - K_{\gamma_{k},\tilde{\theta}_{k}}\hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k+1}) ;\\ \tilde{\eta}_{k}^{d} = \pi_{\gamma_{k},\tilde{\theta}_{k}}(H_{\tilde{\theta}_{k}}) - \pi_{\tilde{\theta}_{k}}(H_{\tilde{\theta}_{k}}) . \end{cases}$$
(34)

Note that if a.s. the following sums converge

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$$\sum_{k=0}^{n} \delta_{k+1} \langle \Pi_{\Theta} \left[\tilde{\theta}_{k} - \delta_{k} \nabla f(\tilde{\theta}_{k}) \right], \tilde{\eta}_{k}^{i} \rangle , \qquad \sum_{k=0}^{n} \delta_{k+1} \tilde{\eta}_{k}^{i} ,$$

where $i \in \{a, b, c, d\}$, then combining this result with (33) and (29) shows that (28) holds which proves the theorem.

(a) By definition (23) and (34), we have that $(\tilde{\eta}_k^a)_{k\in\mathbb{N}}$ is a martingale increment with respect to $(\tilde{\mathcal{F}}_j)_{j\in\mathbb{N}}$ defined by (25). Then, using (32), A4 and H3, we obtain for any $n \in \mathbb{N}$

$$\mathbb{E}\left[\left\|\sum_{k=0}^{n} \delta_{k+1} \tilde{\eta}_{k}^{a}\right\|^{2}\right] = \sum_{k=0}^{n} \delta_{k+1}^{2} \mathbb{E}\left[\|\tilde{\eta}_{k}^{a}\|^{2}\right]$$
$$\leqslant \sum_{k=0}^{n} \delta_{k+1}^{2} \mathbb{E}\left[\mathbb{E}\left[\|\hat{H}_{\tilde{\theta}_{k}}\|^{2} (\tilde{X}_{k+1}) \middle| \mathcal{F}_{k}\right]\right] \leqslant C_{a} \sum_{k=0}^{n} \delta_{k+1}^{2} \gamma_{k}^{-2}$$

,

with $C_a = C_P A_5 \mathbb{E}[V(\tilde{X}_0)]$. The martingale $(\sum_{k=0}^n \delta_{k+1} \tilde{\eta}_k^a)_{n \in \mathbb{N}}$ is then bounded in $L^2(\Omega, \mathbb{P})$ and thus a.s. converges. Similarly, the martingale $(\sum_{k=0}^n \delta_{k+1} \langle \Pi_{\Theta} \left[\tilde{\theta}_k - \delta_k \nabla f(\tilde{\theta}_k) \right], \tilde{\eta}_k^a \rangle)_{n \in \mathbb{N}}$ is also bounded in $L^2(\Omega, \mathbb{P})$ and converges a.s..

(b) Regarding $(\tilde{\eta}_k^b)_{k \in \mathbb{N}}$, using (34), and a discrete integration formula, we obtain for any $n \in \mathbb{N}$,

$$\begin{split} \sum_{k=0}^{n} \delta_{k+1} \tilde{\eta}_{k}^{b} &= \sum_{k=0}^{n} \delta_{k+1} \left(K_{\gamma_{k},\tilde{\theta}_{k}} \hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k}) - K_{\gamma_{k+1},\tilde{\theta}_{k+1}} \hat{H}_{\tilde{\theta}_{k+1}}(\tilde{X}_{k+1}) \right) \\ &= \sum_{k=1}^{n} \left(\delta_{k+1} - \delta_{k} \right) K_{\gamma_{k},\tilde{\theta}_{k}} \hat{H}_{\tilde{\theta}_{k}}(\tilde{X}_{k}) - \delta_{n+1} K_{\gamma_{n+1},\tilde{\theta}_{n+1}} \hat{H}_{\tilde{\theta}_{n+1}}(\tilde{X}_{n+1}) + \delta_{1} K_{\gamma_{0},\tilde{\theta}_{0}} \hat{H}_{\tilde{\theta}_{0}}(\tilde{X}_{0}) \;. \end{split}$$

We have using (32), A4 and H3 for any $n \in \mathbb{N}$

. . .

$$\sum_{k=1}^{n} |\delta_{k+1} - \delta_k| \mathbb{E} \left[K_{\gamma_k, \tilde{\theta}_k} \| \hat{H}_{\tilde{\theta}_k} \| (\tilde{X}_k) \right] \leqslant C_b \sum_{k=1}^{n} |\delta_{k+1} - \delta_k| \gamma_k^{-1} ,$$

where $C_b = C_P A_5 \mathbb{E}[V(\tilde{X}_0)]$. Since by (26) and $\sup_{k \in \mathbb{N}} \gamma_k < \bar{\gamma}$, $\sum_{k=0}^{+\infty} \delta_k^2 / \gamma_k < +\infty$, then the condition $\sup_{k \in \mathbb{N}} |\delta_{k+1} - \delta_k| \delta_k^{-2} < +\infty$ implies

$$\sum_{k=1}^{+\infty} |\delta_{k+1} - \delta_k| \mathbb{E} \left[K_{\gamma_k, \tilde{\theta}_k} \| \hat{H}_{\tilde{\theta}_k} \| (\tilde{X}_k) \right] \leqslant C_b \sup_{k \in \mathbb{N}} \{ |\delta_{k+1} - \delta_k| \, \delta_k^{-2} \} \sum_{k \in \mathbb{N}^*} \delta_k^2 \gamma_k^{-1} < +\infty$$

Therefore, $\sum_{k \in \mathbb{N}^*} (\delta_k - \delta_{k+1}) K_{\gamma_k, \tilde{\theta}_k} \hat{H}_{\tilde{\theta}_k}(\tilde{X}_k)$ is a.s. absolutely convergent. On the other hand, we have that for any $n \in \mathbb{N}$, $K_{\gamma_{n+1}, \tilde{\theta}_{n+1}} \| \hat{H}_{\tilde{\theta}_{n+1}} \| (\tilde{X}_{n+1})$ is bounded by $\gamma_{n+1}^{-1} M_H K_{\gamma_{n+1}, \tilde{\theta}_{n+1}} V^{1/2}(\tilde{X}_{n+1})$. It follows that for any $\varepsilon > 0$, using the Markov inequality and **H3**, we get

$$\begin{split} \sum_{n\in\mathbb{N}} \mathbb{P}\left(\delta_n \gamma_n^{-1} K_{\gamma_{n+1},\tilde{\theta}_{n+1}} \| \hat{H}_{\tilde{\theta}_{n+1}} \| (\tilde{X}_{n+1}) \geqslant \varepsilon\right) &\leq \sum_{n\in\mathbb{N}} \mathbb{P}\left(\delta_n \gamma_n^{-1} K_{\gamma_{n+1},\tilde{\theta}_n} V^{1/2}(\tilde{X}_{n+1}) \geqslant \varepsilon\right) \\ &\leq \varepsilon^{-2} \sum_{n\in\mathbb{N}} \delta_n^2 \gamma_n^{-2} M_H A_5 \mathbb{E}\left[V(\tilde{X}_0)\right] < +\infty \;, \end{split}$$

by (26). Using the Borel-Cantelli lemma we get that a.s. $\lim_{n\to+\infty} \delta_n K_{\gamma_n,\tilde{\theta}_n} \hat{H}_{\tilde{\theta}_n}(\tilde{X}_n) = 0$. This completes the proof of the convergence of $\sum_{k\in\mathbb{N}} \delta_{k+1} \tilde{\eta}_k^b$.

We conclude similarly for the convergence of $\sum_{k\in\mathbb{N}} \delta_{k+1} \langle \Pi_{\Theta} \left[\tilde{\theta}_k - \delta_{k+1} \nabla f(\tilde{\theta}_k) \right], \tilde{\eta}_k^b \rangle$, upon verifying that for any $k \in \mathbb{N}$, $a_{k+1} = \delta_{k+1} \Pi_{\Theta} \left[\vartheta_{k+1} \right]$ where $\vartheta_{k+1} = \tilde{\theta}_k - \delta_{k+1} \nabla f(\tilde{\theta}_k)$ satisfies that there exists $A \ge 0$ such that for any $k \in \mathbb{N}$,

$$||a_{k+1} - a_k|| \leq A\{|\delta_{k+1} - \delta_k|(1+\delta_k) + \delta_k^2\}$$

Indeed, using A1-A2 and that Π_{Θ} is non-expansive by A1 and [23, Proposition 4.8], we get for any $k \in \mathbb{N}$,

$$\begin{aligned} \|a_{k+1} - a_k\| &\leq |\delta_{k+1} - \delta_k| \left\| \Pi_{\Theta}(\vartheta_{k+1}) \right\| + \delta_k \left\| \Pi_{\Theta}(\vartheta_{k+1}) - \Pi_{\Theta}(\tilde{\theta}_k) \right\| \\ &\leq M_{\Theta} \left| \delta_{k+1} - \delta_k \right| + \delta_k \left\| \vartheta_{k+1} - \tilde{\theta}_k \right\| \leq M_{\Theta} \left| \delta_{k+1} - \delta_k \right| + \delta_k \delta_{k+1} \sup_{\Theta} \left\| \nabla f \right\| . \end{aligned}$$

(c) We now turn to $(\tilde{\eta}_k^c)_{k\in\mathbb{N}}$. We start by giving an upper-bound on $\|\pi_{\gamma_1,\theta_1} - \pi_{\gamma_2,\theta_2}\|_{V^{1/2}}$ for $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2$ and $\theta_1, \theta_2 \in \Theta$. Let $f : \mathbb{R}^d \to \mathbb{R}$ be a measurable function satisfying $\sup_{x\in\mathbb{R}^d}\{|f|/V^{1/2}\} \leq 1$. Using H5, H4, H3 and Lemma 11, we get that for any $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2, \theta_1, \theta_2 \in \Theta, x \in \mathbb{R}^d$ and $\ell \in \mathbb{N}^*$

$$\begin{split} \left| K_{\gamma_{1},\theta_{1}}^{\ell}f(x) - K_{\gamma_{2},\theta_{2}}^{\ell}f(x) \right| &= \left| \sum_{j=0}^{\ell} K_{\gamma_{1},\theta_{1}}^{j} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) \left\{ K_{\gamma_{2},\theta_{2}}^{(\ell-j)}f(x) - \pi_{\gamma_{2},\theta_{2}}(f) \right\} \right| \\ &\leq A_{7} \sum_{j=0}^{\ell} \kappa^{(\ell-j)\gamma_{2}} \left| K_{\gamma_{1},\theta_{1}}^{j} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) V^{1/2}(x) \right| \\ &\leq A_{7}A_{6} \sum_{j=0}^{\ell} \kappa^{(\ell-j)\gamma_{2}} \left[\Lambda_{1}(\gamma_{1},\gamma_{2}) + \Lambda_{2}(\gamma_{1},\gamma_{2}) \|\theta_{1} - \theta_{2}\| \right] \sup_{k \in \mathbb{N}} K_{\gamma_{1},\theta_{1}}^{k} V(x) \\ &\leq A_{7}A_{6}A_{5} \log(1/\kappa)^{-1} \kappa^{-\bar{\gamma}} \gamma_{2}^{-1} \left[\Lambda_{1}(\gamma_{1},\gamma_{2}) + \Lambda_{2}(\gamma_{1},\gamma_{2}) \|\theta_{1} - \theta_{2}\| \right] V(x) \,. \end{split}$$

Taking $\ell \to +\infty$ and using H5, we obtain that for any $\theta_1, \theta_2 \in \Theta$ and $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2$,

$$\|\pi_{\gamma_1,\theta_1} - \pi_{\gamma_2,\theta_2}\|_V \leqslant C_{\pi} \gamma_2^{-1} \left[\Lambda_1(\gamma_1,\gamma_2) + \Lambda_2(\gamma_1,\gamma_2) \|\theta_1 - \theta_2\|\right] ,$$
(35)

with $C_{\pi} = A_7 A_6 A_5 \log(1/\kappa)^{-1} \kappa^{-\bar{\gamma}} V(0)$. By (31), (32) and **H** 5, we have for any $\theta_1, \theta_2 \in \Theta$, $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2$ and $x \in \mathbb{R}^d$,

$$\begin{split} \left\| K_{\gamma_{1},\theta_{1}} \hat{H}_{\gamma,\theta_{1}}(x) - K_{\gamma_{2},\theta_{2}} \hat{H}_{\gamma_{2},\theta_{2}}(x) \right\| \\ &= \left\| \sum_{\ell \in \mathbb{N}^{*}} \left\{ K_{\gamma_{1},\theta_{1}}^{\ell} H_{\theta_{1}}(x) - \pi_{\gamma_{1},\theta_{1}}(H_{\theta_{1}}) \right\} - \sum_{\ell \in \mathbb{N}^{*}} \left\{ K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{2}}(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{2}}) \right\} \right\| \\ &\leqslant \sum_{\ell \in \mathbb{N}^{*}} \left\| \left\{ K_{\gamma_{1},\theta_{1}}^{\ell} H_{\theta_{1}}(x) - \pi_{\gamma_{1},\theta_{1}}(H_{\theta_{1}}) \right\} - \left\{ K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{2}}(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{2}}) \right\} \right\| . \end{split}$$

We bound now each term of the sum in the right hand side. Note first that for any bounded measurable functions f_1, f_2 from \mathbb{R}^d to $\mathbb{R}, \theta_1, \theta_2 \in \Theta, \gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2, x \in \mathbb{R}^d$ and

 $\ell \in \mathbb{N}^*,$ it holds that

$$\begin{split} K^{\ell}_{\gamma_{1},\theta_{1}}f_{1}(x) - K^{\ell}_{\gamma_{2},\theta_{2}}f_{2}(x) &= K^{\ell}_{\gamma_{1},\theta_{1}}f_{1}(x) - K^{\ell}_{\gamma_{2},\theta_{2}}f_{1}(x) + K^{\ell}_{\gamma_{2},\theta_{2}}(f_{1}(x) - f_{2}(x)) \\ &= \sum_{j=0}^{\ell-1} \left\{ K^{j}_{\gamma_{1},\theta_{1}} - \pi_{\gamma_{1},\theta_{1}} \right\} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) \left\{ K^{\ell-1-j}_{\gamma_{2},\theta_{2}}f_{1}(x) - \pi_{\gamma_{2},\theta_{2}}(f_{1}) \right\} \\ &\quad + \sum_{j=0}^{\ell-1} \pi_{\gamma_{1},\theta_{1}} \left\{ K^{\ell-1-j}_{\gamma_{2},\theta_{2}}f_{1}(x) - K^{\ell-j}_{\gamma_{2},\theta_{2}}f_{1}(x) \right\} + K^{\ell}_{\gamma_{2},\theta_{2}}(f_{1}(x) - f_{2}(x)) \\ &= \sum_{j=0}^{\ell-1} \left\{ K^{j}_{\gamma_{1},\theta_{1}} - \pi_{\gamma_{1},\theta_{1}} \right\} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) \left\{ K^{\ell-1-j}_{\gamma_{2},\theta_{2}}f_{1}(x) - \pi_{\gamma_{2},\theta_{2}}(f_{1}) \\ &- \pi_{\gamma_{1},\theta_{1}}(K^{\ell}_{\gamma_{2},\theta_{2}}f_{1}(x) - f_{1}(x)) + K^{\ell}_{\gamma_{2},\theta_{2}}(f_{1}(x) - f_{2}(x)) \,. \end{split}$$

Setting $\tilde{H}_{\theta_1} = H_{\theta_1} - \pi_{\gamma_1,\theta_1}(H_{\theta_1})$ and $\tilde{H}_{\theta_2} = H_{\theta_2} - \pi_{\gamma_2,\theta_2}(H_{\theta_2})$, we obtain that

$$K_{\gamma_{1},\theta_{1}}^{\ell}\tilde{H}_{\theta_{1}}(x) - K_{\gamma_{2},\theta_{2}}^{\ell}\tilde{H}_{\theta_{2}}(x) = \sum_{j=0}^{\ell-1} \left\{ K_{\gamma_{1},\theta_{1}}^{j} - \pi_{\gamma_{1},\theta_{1}} \right\} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) \left\{ K_{\gamma_{2},\theta_{2}}^{\ell-1-j}H_{\theta_{1}}(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}}) \right\} + \Xi , \quad (36)$$

where

$$\begin{aligned} \Xi &= -\pi_{\gamma_{1},\theta_{1}} (K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{1}}(x) - H_{\theta_{1}}(x)) + K_{\gamma_{2},\theta_{2}}^{\ell} \left[H_{\theta_{1}}(x) - H_{\theta_{2}}(x) + \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{2}}) - \pi_{\gamma_{1},\theta_{1}}(H_{\theta_{1}}) \right] \\ &= -\pi_{\gamma_{1},\theta_{1}} K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{1}}(x) + K_{\gamma_{2},\theta_{2}}^{\ell} \left[H_{\theta_{1}}(x) - H_{\theta_{2}}(x) + \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{2}}) \right] \\ &= (\pi_{\gamma_{2},\theta_{2}} - \pi_{\gamma_{1},\theta_{1}}) (K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{1}}(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}})) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}}) \\ &\quad + K_{\gamma_{2},\theta_{2}}^{\ell} \left[H_{\theta_{1}}(x) - H_{\theta_{2}}(x) + \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{2}}) \right] \\ &= (\pi_{\gamma_{2},\theta_{2}} - \pi_{\gamma_{1},\theta_{1}}) (K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{1}}(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}})) + K_{\gamma_{2},\theta_{2}}^{\ell} (H_{\theta_{1}} - H_{\theta_{2}})(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}} - H_{\theta_{2}}). \end{aligned}$$

$$(37)$$

For the first term in the decomposition (36), using A4, H5 and H4, we obtain for any $\theta_1, \theta_2 \in \Theta$, $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2$, $x \in \mathbb{R}^d$ and $\ell \in \mathbb{N}^*$

$$\left\| \sum_{j=0}^{\ell-1} \left\{ K_{\gamma_{1},\theta_{1}}^{j} - \pi_{\gamma_{1},\theta_{1}} \right\} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) \left\{ K_{\gamma_{2},\theta_{2}}^{\ell-1-j} H_{\theta_{1}}(x) - \pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}}) \right\} \right\|$$

$$\leq A_{7} M_{H} \sum_{j=0}^{\ell-1} \left| \left\{ K_{\gamma_{1},\theta_{1}}^{j} - \pi_{\gamma_{1},\theta_{1}} \right\} (K_{\gamma_{1},\theta_{1}} - K_{\gamma_{2},\theta_{2}}) V^{1/4}(x) \right|$$

$$\leq A_{7}^{2} A_{6} M_{H} \left[\Lambda_{1}(\gamma_{1},\gamma_{2}) + \Lambda_{2}(\gamma_{1},\gamma_{2}) \|\theta_{1} - \theta_{2}\| \right] \sum_{j=0}^{\ell-1} \kappa^{j\gamma} \kappa^{(\ell-1-j)\gamma_{2}} V^{1/2}(x)$$

$$\leq A_{7}^{2} A_{6} M_{H} \left[\Lambda_{1}(\gamma_{1},\gamma_{2}) + \Lambda_{2}(\gamma_{1},\gamma_{2}) \|\theta_{1} - \theta_{2}\| \right] \ell \kappa^{\ell\gamma_{2}} V^{1/2}(x) . \tag{38}$$

For the first term in (37), using A4, H5 and (35) we obtain for any $\theta_1, \theta_2 \in \Theta, \gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2, x \in \mathbb{R}^d$ and $\ell \in \mathbb{N}^*$

$$\left\| (\pi_{\gamma_{1},\theta_{1}} - \pi_{\gamma_{2},\theta_{2}}) (K_{\gamma_{2},\theta_{2}}^{\ell} H_{\theta_{1}}(x) - \pi_{\gamma_{2},\theta_{2}} (H_{\theta_{1}})) \right\| \leq A_{7} M_{H} \kappa^{\ell\gamma_{2}} \|\pi_{\gamma_{1},\theta_{1}} - \pi_{\gamma_{2},\theta_{2}}\|_{V^{1/2}} \leq A_{7} M_{H} C_{\pi} \kappa^{\ell\gamma_{2}} \gamma_{2}^{-1} \left[\Lambda_{1}(\gamma_{1},\gamma_{2}) + \Lambda_{2}(\gamma_{1},\gamma_{2}) \|\theta_{1} - \theta_{2} \| \right] .$$

$$(39)$$

For the second term in (37), using A4 and H5, we obtain for any $\theta_1, \theta_2 \in \Theta, \gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2, x \in \mathbb{R}^d$ and $\ell \in \mathbb{N}^*$

$$\left\|K_{\gamma_{2},\theta_{2}}^{\ell}(H_{\theta_{1}}-H_{\theta_{2}})(x)-\pi_{\gamma_{2},\theta_{2}}(H_{\theta_{1}}-H_{\theta_{2}})\right\| \leqslant A_{7}L_{H}\kappa^{\ell\gamma_{2}}\|\theta_{1}-\theta_{2}\|V^{1/2}(x).$$
(40)

Combining (38), (39), (40), (37) in (36), and using Lemma 11, we obtain that for any $\theta_1, \theta_2 \in \Theta$, $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2, x \in \mathbb{R}^d$ that

$$\begin{split} \left\| K_{\gamma_1,\theta_1} \hat{H}_{\gamma_1,\theta_1}(x) - K_{\gamma_2,\theta_2} \hat{H}_{\gamma_2,\theta_2}(x) \right\| \\ & \leq C_c \gamma_2^{-1} \left[\gamma_2^{-1} \left\{ \Lambda_1(\gamma_1,\gamma_2) + \Lambda_2(\gamma_1,\gamma_2) \| \theta_1 - \theta_2 \| \right\} + \| \theta_1 - \theta_2 \| \right] V^{1/2}(x) , \end{split}$$

with

$$C_c = A_7 (L_H + C_\pi M_H) \log(1/\kappa)^{-1} \kappa^{-\bar{\gamma}} + A_7^2 A_6 M_H \log(1/\kappa)^{-2} \kappa^{-\bar{\gamma}} .$$

Combining this result and that for any $k \in \mathbb{N}$, $\|\tilde{\theta}_{k+1} - \tilde{\theta}_k\| \leq \delta_{k+1} M_H V^{1/2}(\tilde{X}_{k+1})$ using (9), A4 and Π_{Θ} is non-expansive, we get that for any $k \in \mathbb{N}$,

$$\begin{split} \left\| K_{\gamma_{k},\tilde{\theta}_{k}} \hat{H}_{\gamma_{k},\tilde{\theta}_{k}}(x) - K_{\gamma_{k+1},\tilde{\theta}_{k+1}} \hat{H}_{\gamma_{k+1},\tilde{\theta}_{k+1}}(x) \right\| \\ & \leq C_{c}(1+M_{H})\gamma_{k+1}^{-1} \left[\gamma_{k+1}^{-1} \left\{ \Lambda_{1}(\gamma_{k},\gamma_{k+1}) + \Lambda_{2}(\gamma_{k},\gamma_{k+1})\delta_{k+1} \right\} + \delta_{k+1} \right] V(\tilde{X}_{k+1}) \; . \end{split}$$

Then since $\sup_{n \in \mathbb{N}} \mathbb{E}\left[V(\tilde{X}_n)\right] < +\infty$ by **H3**, (34) and (26), we get

$$\sum_{k\in\mathbb{N}}\delta_{k+1}\mathbb{E}\left[\|\tilde{\eta}_k^c\|\right] < +\infty.$$

(d) By a straightforward application of H1-(ii), H5 and (26), we conclude that $\sum_{k \in \mathbb{N}} \delta_{k+1} \| \tilde{\eta}_k^d \|$ converges.

B.2 Proofs of Theorem 5

In this section, we give the proof of Theorem 5 by showing that H1 and H2 hold. However, first of all, we establish stability results uniform in the parameter $\theta \in \Theta$ for Langevin diffusion (5) and the associated Euler-Maruyama discretization (6) based on Foster-Lyapunov drift condition with constants independent of θ .

Under L 1, for any $\theta \in \Theta$, (5) defines a Markov semi-group $(P_{t,\theta})_{t\geq 0}$ for any $x \in \mathbb{R}^d$ and $\mathsf{A} \in \mathcal{B}(\mathbb{R}^d)$ by $P_{t,\theta}(x,\mathsf{A}) = \mathbb{P}(Y_t^{\theta} \in \mathsf{A})$ where $(Y_t^{\theta})_{t\geq 0}$ is the solution of (5) with $Y_0^{\theta} = x$. Consider now the generator of $(P_{t,\theta})_{t\geq 0}$ for any $\theta \in \Theta$, defined for any $f \in C^2(\mathbb{R}^d)$ by

$$\mathcal{A}_{\theta}f = \langle \nabla_x f, \nabla_x U_{\theta}(x) \rangle + \Delta_x f .$$

We say that a Markov kernel R on $\mathbb{R}^d \times \mathcal{B}(\mathbb{R}^d)$ satisfies a discrete Foster-Lyapunov drift condition $\mathbf{D}_{\mathrm{d}}(V,\lambda,b)$ if there exit $\lambda \in (0,1), b \ge 0$ and a measurable $V : \mathbb{R}^d \to [1,+\infty)$ such that for all $x \in \mathbb{R}^d$

$$RV(x) \leq \lambda V(x) + b$$
.

We say that a Markov semi-group $(P_t)_{t\geq 0}$ on $\mathbb{R}^d \times \mathcal{B}(\mathbb{R}^d)$ with extended infinitesimal generator $(\mathcal{A}, \mathcal{D}(\mathcal{A}))$ (see e.g. [32] for the definition of $(\mathcal{A}, \mathcal{D}(\mathcal{A}))$) satisfies a continuous drift condition $\mathbf{D}_c(V, \zeta, \beta)$ if there exist $\zeta > 0, \beta \geq 0$ and a measurable function $V : \mathbb{R}^d \to [1, +\infty)$ with $V \in \mathcal{D}(\mathcal{A})$ such that for all $x \in \mathbb{R}^d$

$$\mathcal{A}V(x) \leq -\zeta V(x) + \beta$$
.

B.2.1 Foster-Lyapunov drift conditions uniform on θ

Proposition 12. Assume L2. Let $\bar{\gamma} < \min(1, 2m)$ and define $V : \mathbb{R}^d \to [1, +\infty)$ such that for all $x \in \mathbb{R}^d$, $V(x) = \exp(\tilde{\kappa}\phi(x)/4)$ with $\phi(x) = \sqrt{\|x\|^2 + 1}$ and $\tilde{\kappa} = \eta/4$. Then there exist $\lambda \in (0, 1)$ and $b \ge 0$ such that for all $\gamma \in (0, \bar{\gamma})$ and $\theta \in \Theta$ the Markov kernel $R_{\gamma,\theta}$ satisfies the discrete drift condition $\mathbf{D}_{d}(V, \lambda^{\gamma}, b\gamma)$, i.e. for all $x \in \mathbb{R}^d$

$$R_{\gamma,\theta}V(x) \leq \lambda^{\gamma}V(x) + b\gamma$$
,

with

$$\lambda = e^{-2^{-4}\eta^2(2^{1/2}-1)} , \quad b = \tilde{\kappa}(d+c+2^{1/2}\tilde{\kappa}) \exp\left[\tilde{\kappa}\left\{(d+c+\tilde{\kappa})\bar{\gamma} + \phi(\max(1,2(d+c)/\eta,M_\eta))\right\}\right]$$

Proof. Since ϕ is 1-Lipschitz, by the log-Sobolev inequality [4, Proposition 5.4.1] and the concavity of $x \mapsto x^{1/2}$, we have for any $x \in \mathbb{R}^d$

$$R_{\gamma,\theta}V(x) \leqslant e^{\tilde{\kappa}R_{\gamma,\theta}\phi(x)+\tilde{\kappa}^2\gamma} \leqslant e^{\tilde{\kappa}\sqrt{\|x-\gamma\nabla_x U_\theta(x)\|^2+2\gamma d+1}+\tilde{\kappa}^2\gamma}.$$

Using **L2** and $\gamma < 2m$ we obtain that for any $x \in \mathbb{R}^d$,

$$\begin{aligned} \|x - \gamma \nabla_x U_{\theta}(x)\|^2 &\leq \|x\|^2 - 2\gamma \langle x, \nabla_x U_{\theta}(x) \rangle + \gamma^2 \|\nabla_x U_{\theta}(x)\|^2 \\ &\leq \|x\|^2 - 2\eta \gamma \|x\| \mathbb{1}_{\|x\| \ge M_{\eta}} + \gamma(\gamma - 2m) \|\nabla_x U_{\theta}(x)\|^2 + 2\gamma c \\ &\leq \|x\|^2 - 2\eta \gamma \|x\| \mathbb{1}_{\|x\| \ge M_{\eta}} + 2\gamma c . \end{aligned}$$

Therefore, for any $x \in \mathbb{R}^d$, $||x|| \ge \max(1, 2(d+c)/\eta, M_\eta)$, using for any a > 0, $\sqrt{1+a} - 1 \le a/2$, we get for any $\theta \in \Theta$

$$\sqrt{\|x - \gamma \nabla_x U_{\theta}(x)\|^2 + 2\gamma d + 1} - \phi(x)
\leq \phi(x) \left\{ \sqrt{\phi^{-2}(x) (\|x\|^2 + 1 - 2\eta\gamma\|x\| + 2\gamma(d+c))} - 1 \right\}
\leq \phi(x) \left\{ \sqrt{1 + 2\gamma \phi^{-2}(x)(d+c-\eta\|x\|)} - 1 \right\}
\leq \gamma \phi^{-1}(x)(d+c-\eta\|x\|) \leq -\eta\gamma \phi^{-1}(x)\|x\|/2 \leq -2^{-3/2}\eta\gamma.$$
(41)

This gives $R_{\gamma,\theta}V(x) \leq e^{-2^{-4}\eta^2(2^{1/2}-1)\gamma} = \lambda^{\gamma}$. For all $x \in \mathbb{R}^d$ we obtain with the same arguments

$$\sqrt{\|x - \gamma \nabla_x U_\theta(x)\|^2 + 2\gamma d + 1} - \phi(x) \leq \gamma(d+c) , \qquad (42)$$

and consequently, combining (42) and (41)

$$R_{\gamma,\theta}V(x) \leqslant \lambda^{\gamma}V(x) + \left(e^{\tilde{\kappa}(d+c+\tilde{\kappa})\gamma} - \lambda^{\gamma}\right)e^{\tilde{\kappa}\phi(\max(1,2(d+c)/\eta,M_{\eta}))}\mathbb{1}_{\|x\| \leqslant \max(1,2(d+c)/\eta,M_{\eta})} \cdot$$

Using $e^t - 1 \leq te^t$ for $t \geq 0$ we obtain the result.

Proposition 13. Let $\bar{\gamma} < \min(1, \eta/L^2)$ and define $V : \mathbb{R}^d \to [1, +\infty)$ such that for all $x \in \mathbb{R}^d$, $V(x) = 1 + ||x||^2$. Assume **L1** and **L3**. Then there exist $\lambda \in (0, 1)$ and $b \ge 0$ such that for all $\gamma \in (0, \bar{\gamma})$ and $\theta \in \Theta$ the Markov kernel $R_{\gamma,\theta}$ satisfies the discrete drift condition $\mathbf{D}_{d}(V, \lambda^{\gamma}, b\gamma)$, i.e. for all $x \in \mathbb{R}^d$

$$R_{\gamma,\theta}V(x) \leqslant \lambda^{\gamma}V(x) + b\gamma$$

with $\lambda = \exp(-2(\eta - \bar{\gamma}L^2))$ and $b = 2(c + d + \eta(M_{\eta}^2 + 1) + \bar{\gamma} \sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\|^2).$

Proof. Set Z a standard Gaussian random variable over \mathbb{R}^d we have using L1 and $(a+b)^2 \leq 2(a^2+b^2)$

$$\begin{aligned} R_{\gamma,\theta}V(x) &= 1 + \mathbb{E}(\|x - \gamma \nabla_x U_{\theta}(x) + \sqrt{2\gamma Z}\|^2) \\ &\leq 1 + \|x\|^2 - 2\gamma \langle x, \nabla_x U_{\theta}(x) \rangle + 2\gamma d + \gamma^2 \|\nabla_x U_{\theta}(x)\|^2 \\ &\leq V(x) + 2\gamma c - 2\gamma \eta \|x\|^2 \mathbb{1}_{\overline{\mathrm{B}}(0,M_{\eta})^c}(x) + 2\gamma d + 2\gamma^2 \mathrm{L}^2 \|x\|^2 + 2\gamma^2 \sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\|^2 \\ &\leq (1 - 2\gamma (\eta - \gamma \mathrm{L}^2)) V(x) + 2\gamma (c + d + \eta (M_{\eta}^2 + 1) + \gamma \sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\|^2) , \end{aligned}$$

where Z is a zero mean d-dimensional Gaussian random variable with covariance identity. \Box

Proposition 14. Let $\bar{\gamma} < \min(1, \eta/L^2)$ and define $V : \mathbb{R}^d \to [1, +\infty)$ such that for all $x \in \mathbb{R}^d$, $V(x) = 1 + ||x||^4$. Assume **L1** and **L3**. Then there exist $\lambda \in (0, 1)$ and $b \ge 0$ such that for all $\gamma \in (0, \bar{\gamma})$ and $\theta \in \Theta$ the Markov kernel $R_{\gamma,\theta}$ satisfies the discrete drift condition $\mathbf{D}_d(V, \lambda^{\gamma}, b\gamma)$, i.e. for all $x \in \mathbb{R}^d$

$$R_{\gamma,\theta}V(x) \leqslant \lambda^{\gamma}V(x) + b\gamma$$
,

Proof. Let $a, b \in \mathbb{R}^d$. We have the following equality

$$||a+b||^{2} = ||a||^{4} + ||b||^{4} + 4\langle a,b\rangle^{2} + 2||a||^{2}||b||^{2} + 4\langle a,b\rangle||a||^{2} + 4\langle a,b\rangle||b||^{2}.$$
(43)

Let a = x and $b = -\gamma \nabla_x U_{\theta}(x) + \sqrt{2\gamma}Z$ where Z is a zero mean d-dimensional Gaussian random variable with covariance identity. Using (43) we get for any $\theta \in \Theta$ and $\gamma > 0$

$$\begin{split} R_{\gamma,\theta}V(x) &= 1 + \|x\|^4 + \mathbb{E}\left[\|-\gamma\nabla_x U_\theta(x) + \sqrt{2\gamma}Z\|^4\right] + 4\mathbb{E}\left[\langle x, -\gamma\nabla_x U_\theta(x) + \sqrt{2\gamma}Z\rangle^2\right] \\ &+ 2\|x\|^2\mathbb{E}\left[\|-\gamma\nabla_x U_\theta(x) + \sqrt{2\gamma}Z\|^2\right] + 4\mathbb{E}\left[\langle x, -\gamma\nabla_x U_\theta(x) + \sqrt{2\gamma}Z\rangle\| - \gamma\nabla_x U_\theta(x) + \sqrt{2\gamma}Z\|^2\right] \\ &+ 4\mathbb{E}\left[\langle x, -\gamma\nabla_x U_\theta(x) + \sqrt{2\gamma}Z\rangle\|x\|^2\right] \,. \end{split}$$

Let $n(x) = \|\nabla_x U_\theta(x)\|$ and $r(x) = \langle \nabla_x U_\theta(x), x \rangle$. Using the Isserlis formula [25] we obtain that

$$R_{\gamma,\theta}V(x) = V(x) + \gamma^4 n^4(x) + 4\gamma^2 d(d+2) + 4\gamma^2 n^2(x) + 4\gamma^2 n^2(x)d + \gamma^2 r^2(x) + 2\gamma ||x||^2 - 4\gamma r(x) ||x||^2 - 4\gamma^3 n^2(x)r(x) - 8\gamma^2 dr(x) - 16\gamma^2 r(x) + 2\gamma^2 ||x||^2 n^2(x) + 4\gamma d||x||^2 .$$
(44)

Let $||x|| \ge M_{\eta}$ then $r(x) \ge ||x||^2 - c$, $n(x) \le L||x|| + \sup_{\theta \in \Theta} ||\nabla_x U_{\theta}(0)||$ and using the Cauchy-Schwarz inequality we have $r(x) \le L||x||^2 + \sup_{\theta \in \Theta} ||\nabla_x U_{\theta}(0)||$. Therefore, using (44), there exists a polynomial P of order strictly less than 4 such that

$$R_{\gamma,\theta}V(x)\leqslant V(x)+\gamma(\gamma^3\mathsf{L}^4+3\gamma\mathsf{L}^2-4\eta)V(x)+\gamma\mathsf{P}(\|x\|)\;.$$

Since $\bar{\gamma} = \min(1, \eta/L^2)$, for all $\gamma \in (0, \bar{\gamma}), \theta \in \Theta$ and $||x|| \ge M_{\eta}$, we have that

$$R_{\gamma,\theta}V(x) \leqslant \tilde{\lambda}^{\gamma}V(x) + \gamma P(V(x)^{1/4})$$

with $\tilde{\lambda} = \exp(\bar{\gamma}^3 L^4 + 3\bar{\gamma}L^2 - 4\eta) \in (0,1)$. Let $1 > \lambda > \tilde{\lambda}$ and b_1 defined by

$$b_1 = \sup_{\gamma \in (0,\bar{\gamma})} \sup_{x \in \mathbb{R}^d} \gamma^{-1} (\tilde{\lambda}^{\gamma} - \lambda^{\gamma}) V(x) + \mathcal{P}(V(x)^{1/4}) < +\infty$$

We obtain that for any $\gamma \in (0, \bar{\gamma}), \theta \in \Theta$ and $||x|| \ge M_{\eta}$ we have

$$R_{\gamma,\theta}V(x) \leq \lambda^{\gamma}V(x) + b_1\gamma$$
.

In addition there exists $b_2 \ge 0$ such that for any $\gamma \in (0, \bar{\gamma}), \theta \in \Theta$ and $||x|| \le M_\eta$ we have

$$R_{\gamma,\theta}V(x) \leqslant b_2\gamma$$

Let $b = \max(b_1, b_2)$ we obtain that for any $\gamma \in (0, \overline{\gamma}), \ \theta \in \Theta$ and $x \in \mathbb{R}^d$

$$R_{\gamma,\theta}V(x) \leq \lambda^{\gamma}V(x) + b\gamma$$
.

Proposition 15. The following properties hold

- 1. Assume L2 then there exist $\zeta > 0$ and $\beta \ge 0$ such that for all $\theta \in \Theta$, $(P_{t,\theta})_{t>0}$ satisfies the continuous drift condition $\mathscr{D}(V,\zeta,\beta)$ for V defined in Proposition 12.
- 2. Assume L1 and L3 then there exist $\zeta > 0$ and $\beta \ge 0$ such that for all $\theta \in \Theta$, $(P_{t,\theta})_{t>0}$ satisfies the continuous drift condition $\mathscr{D}(V,\zeta,\beta)$ for V defined in Proposition 13. In addition ζ and β are defined by

$$\zeta = 2\eta$$
, $\beta = 2(d + c + \eta(M_{\eta}^2 + 1))$.

3. Assume L1 and L3 then there exist $\zeta > 0$ and $\beta \ge 0$ such that for all $\theta \in \Theta$, $(P_{t,\theta})_{t>0}$ satisfies the continuous drift condition $\mathscr{D}(V,\zeta,\beta)$ for V defined in Proposition 14.

Proof. (a) Assuming L 2 the drift condition is a direct consequence of [12, Proposition 16] for $V(x) = \exp(\eta \phi(x)/4)$ as in Proposition 12.

(b) Assuming L3 and setting $V(x) = 1 + ||x||^2$ we get for any $x \in \mathbb{R}^d$ and $\theta \in \Theta$

$$-\langle \nabla_x U_\theta(x), \nabla V(x) \rangle + \Delta_x V(x) = -2\langle \nabla_x U_\theta(x), x \rangle + 2d$$

$$\leqslant -2\eta \mathbb{1}_{\overline{B}(0,M_\eta)^c} V(x) + 2\eta + 2d + 2c$$

$$\leqslant -2\eta V(x) + 2(d+c+\eta(M_\eta^2+1)) .$$

(c) Assuming L3 and setting $V(x) = 1 + ||x||^4$ there exists a polynom P of order strictly less than 4 such that for any $x \in \mathbb{R}^d$ and $\theta \in \Theta$

$$\begin{aligned} -\langle \nabla_x U_\theta(x), \nabla V(x) \rangle + \Delta_x V(x) &= -4 \|x\|^2 \langle \nabla_x U_\theta(x), x \rangle + 4 \|x\|^2 \sum_{i=1}^d x_i \\ &\leqslant -4\eta \mathbb{1}_{\overline{B}(0, M_\eta)^c} V(x) + 4 \|x\|^2 \sum_{i=1}^d x_i + 4c + 4\eta \|x\|^2 \\ &\leqslant -4\eta V(x) + \mathcal{P}(V(x)^{1/4}) . \end{aligned}$$

Since P is of order strictly less than 4 there exists $B \ge$ such that for all $x \in \mathbb{R}^d$ and $\theta \in \Theta$

$$-\langle \nabla_x U_\theta(x), \nabla V(x) \rangle + \Delta_x V(x) \leqslant -3\eta V(x) + B.$$

B.2.2 Checking H1

Lemma 16. Let $V : \mathbb{R}^d \to [1, +\infty)$ satisfying $\lim_{\|x\|\to+\infty} V(x) = +\infty$ and $V \in D(\mathcal{A})$.

(a) Assume there exist $\lambda \in (0,1)$, $b \ge 0$ and $\bar{\gamma} > 0$ such that for any $\theta \in \Theta$ and $\gamma \in (0,\bar{\gamma})$, $R_{\gamma,\theta}$ satisifies $\mathbf{D}_{\mathrm{d}}(V,\lambda^{\gamma},b\gamma)$. Then for any $\theta \in \Theta$ and $\gamma \in (0,\bar{\gamma})$, $R_{\gamma,\theta}$ admits an invariant probability measure $\pi_{\gamma,\theta}$ on $(\mathbb{R}^d, \mathcal{B}(\mathbb{R}^d))$ and there exists $A_0 \ge 0$ such that for any $x \in \mathbb{R}^d$ and $k \in \mathbb{N}$

$$\delta_x R^k_{\gamma,\theta} V \leqslant A_0 + V(x) , \qquad \pi_{\gamma,\theta}(V) \leqslant A_0 , \qquad A_0 = b\lambda^{-\bar{\gamma}}/\log(1/\lambda) .$$

In addition, for all $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $\lim_{k \to +\infty} \|\delta_x R_{\gamma,\theta}^k - \pi_{\gamma,\theta}\|_V = 0$.

(b) Assume there exist $\zeta > 0$ and $\beta \ge 0$ such that for any $\theta \in \Theta$, \mathcal{A}_{θ} satisfies $\mathscr{D}(V, \zeta, \beta)$. Then for any $\theta \in \Theta$, the diffusion is non-explosive, \mathcal{A}_{θ} admits π_{θ} as an invariant probability measure and

$$\tau_{\theta}(V) \leqslant \beta/\zeta$$

In addition, for all $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $\lim_{t \to +\infty} \|\delta_x P_{\theta,t} - \pi_{\theta}\|_V = 0$.

Proof. (a) for any $\gamma \in (0, \bar{\gamma})$ and $\theta \in \Theta$, $R_{\gamma,\theta}$ is irreducible with respect to the Lebesgue measure on \mathbb{R}^d , has the Feller property and satisfies $\mathbf{D}_d(V, \lambda^{\gamma}, b\gamma)$ then [30, Section 4.4] applies and $R_{\gamma,\theta}$ admits an invariant probability measure $\pi_{\gamma,\theta}$. The discrete drift condition and [12, Lemma 1] give that for any $\gamma \in (0, \bar{\gamma})$ and $\theta \in \Theta$

$$R^k_{\gamma,\theta}V(x) \leqslant V(x) + b\lambda^{-\bar{\gamma}}/\log(1/\lambda) , \qquad \pi_{\gamma,\theta}(V) \leqslant b\lambda^{-\bar{\gamma}}/\log(1/\lambda) .$$

We obtain that for all $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $\lim_{k \to +\infty} \|\delta_x P_{t,\theta} - \pi_{\gamma,\theta}\|_V = 0$ using [31, Theorem 16.0.1].

(b) Using $\mathscr{D}(V,\zeta,\beta)$ and [32, Theorem 2.1] we get that the diffusion process is non-explosive and thus $(P_{t,\theta})_{t\geq 0}$ is defined for any $\theta \in \Theta$ and $t \geq 0$. Using [48, Corollary 10.1.4] for any $\theta \in \Theta$, $(P_{t,\theta})_{t\geq 0}$ is strongly Feller continuous, therefore any compact sets is petite for the Markov kernel $P_{h,\theta}$, for any h > 0 and $\theta \in \Theta$, by [31, Theorem 6.0.1]. Using [39, Chapter 7, Proposition 1.5], [19, Chapter 4, Theorem 9.17], and the fact that $\pi_{\theta}(\mathcal{A}_{\theta}f) = 0$ for any $\theta \in \Theta$ and $f \in C_c^2(\mathbb{R}^d)$, we obtain that for any $\theta \in \Theta$, π_{θ} is an invariant measure for $(P_{t,\theta})_{t \ge 0}$. Using $\mathscr{D}(V,\zeta,\beta)$ and [32, Theorem 4.5, Theorem 6.1] we get that for all $\theta \in \Theta$, $\pi_{\theta}(V) \le \beta/\zeta$ and the uniform V-geometric ergodicity for $(P_{t,\theta})_{t \ge 0}$.

As an immediate corollary we obtain that under the conditions of Lemma 16 for any $\theta \in \Theta$, $\gamma \in (0, \bar{\gamma})$ and $k \in \mathbb{N}$, $\pi_{\theta} R^k_{\gamma, \theta} V \leq \beta/\zeta + b\lambda^{-\bar{\gamma}}/\log(1/\lambda)$.

Lemma 17. Let $V : \mathbb{R}^d \to [1, +\infty)$. Assume there exist $\lambda \in (0, 1)$, $b \ge 0$ and $\bar{\gamma} > 0$ such that for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$ $R_{\gamma, \theta}$ satisifies $\mathbf{D}_{d}(V, \lambda^{\gamma}, b\gamma)$ then there exists $B_1 \ge 0$ such that $(X_k^n)_{n \in \mathbb{N}, k \in \{0, \dots, m_n\}}$ given by (9) satisfies for all $n \in \mathbb{N}$ and $k \in \{0, \dots, m_n\}$

$$\mathbb{E}\left[V(X_k^n) \middle| X_0^0\right] \leqslant B_1 V(X_0^0) , \qquad B_1 = 1 + b\lambda^{-\bar{\gamma}} / \log(1/\lambda) .$$

Proof. Let $n \in \mathbb{N}$ and $k \in \{0, \ldots, m_{n+1}\}$. By induction we obtain that

$$\mathbb{E}\left[V(X_{k}^{n+1})\big|\mathcal{F}_{n}\right] = R_{\theta_{n+1},\gamma_{n+1}}^{k}V(X_{0}^{n+1}) \leqslant \lambda^{k\gamma_{n+1}}V(X_{0}^{n+1}) + b\gamma_{n+1}\sum_{i=1}^{k}\lambda^{\gamma_{n+1}(k-i)}.$$

In the same manner we obtain for any $k \in \{0, \ldots, m_0\}$

$$\mathbb{E}\left[V(X_k^0)\big|X_0^0\right] = R_{\theta_0,\gamma_0}^k V(X_0^0) \leqslant \lambda^{k\gamma_0} V(X_0^0) + b\gamma_0 \sum_{i=1}^k \lambda^{\gamma_0(k-i)}$$

Let $\vartheta_{n,k} = \sum_{j=0}^{n-1} m_j + k$, $\vartheta_n = \vartheta_{n,0}$ and $\tilde{\gamma}_i = \sum_{j=0}^{+\infty} \gamma_j \mathbb{1}_{(\vartheta_j,\vartheta_{j+1}]}(i)$. Let also $\Gamma_{p,q} = \sum_{i=p}^{q} \tilde{\gamma}_i$ and $\Gamma_p = \Gamma_{1,p}$. By induction we get for any $n \in \mathbb{N}$ and $k \in \{0, \ldots, m_n\}$

$$\mathbb{E}\left[V(X_k^n)\big|X_0^0\right] \leqslant \lambda^{\Gamma_{\vartheta_{n,k}}}V(X_0^0) + b\sum_{i=1}^{\vartheta_{n,k}} \tilde{\gamma}_i \lambda^{\Gamma_{i+1,\vartheta_{n,k}}}$$

Since $(\tilde{\gamma}_i)_{i\in\mathbb{N}}$ is nonincreasing and for all $t \ge 0$, $1 - \lambda^t \ge -t\lambda^t \log(\lambda)$, we have for all $n \in \mathbb{N}$ and $k \in \{0, \ldots, m_n\}$ that

$$\begin{split} &\sum_{i=1}^{n} \tilde{\gamma}_{i} \lambda^{\Gamma_{i+1,n}} \leqslant \sum_{i=1}^{n} \tilde{\gamma}_{i} \prod_{j=i+1}^{n} (1+\lambda^{\tilde{\gamma}_{1}} \log(\lambda) \tilde{\gamma}_{j}) \\ &\leqslant (-\lambda^{\tilde{\gamma}_{1}} \log(\lambda))^{-1} \sum_{i=1}^{n} \left\{ \prod_{j=i+1}^{n} (1+\lambda^{\tilde{\gamma}_{1}} \log(\lambda) \tilde{\gamma}_{j}) - \prod_{j=i}^{n} (1+\lambda^{\tilde{\gamma}_{1}} \log(\lambda) \tilde{\gamma}_{j}) \right\} \\ &\leqslant (-\lambda^{\tilde{\gamma}_{1}} \log(\lambda))^{-1} . \end{split}$$

Proposition 18. Let $V : \mathbb{R}^d \to [1, +\infty)$ and $M_{V,2} \ge 0$ such that for any $x \in \mathbb{R}^d$, $\sup_{x \in \mathbb{R}^d} (1 + ||x||)^2 / V(x) \le M_{V,2}$. Assume L_1 and there exist $\lambda \in (0, 1)$, $b \ge 0$ and $\overline{\gamma} > 0$ such that for any $\theta \in \Theta$ and $\gamma \in (0, \overline{\gamma})$ $R_{\gamma,\theta}$ satisfies $\mathbf{D}_{\mathbf{d}}(V, \lambda^{\gamma}, b\gamma)$. Assume there exist $\zeta > 0$ and $\beta \ge 0$ such that for any $\theta \in \Theta$, \mathcal{A}_{θ} satisfies $\mathcal{D}(V, \zeta, \beta)$. We obtain that there exists $B_3 \ge 0$ such that for any $\theta \in \Theta$ and $\gamma \in (0, \overline{\gamma})$

$$\|\pi_{\gamma,\theta} - \pi_{\theta}\|_{V^{1/2}} \leq B_3 \gamma^{1/2}$$
.

Proof. Since for all $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$ we have that $R_{\gamma,\theta}$ satisfies $\mathbf{D}_{\mathrm{d}}(V, \lambda, \gamma)$ we also have that $R_{\gamma,\theta}$ satisfies $D(V^{1/2}, \lambda^{1/2}, b\lambda^{-\bar{\gamma}/2}/2)$. Using Lemma 16 we obtain that for any $x \in \mathbb{R}^d$ and $\theta \in \Theta$

$$\lim_{k \to +\infty} \|\pi_{\theta} R^k_{\gamma,\theta} - \pi_{\theta} P_{\gamma k,\theta}\|_{V^{1/2}} = \|\pi_{\gamma,\theta} - \pi_{\theta}\|_{V^{1/2}}.$$

Following the lines of [12, Theorem 10] we write $k = q_{\gamma}m_{\gamma} + r_{\gamma}$ with $m_{\gamma} = \lceil 1/\gamma \rceil$ and $q_{\gamma}, r_{\gamma} \in \mathbb{N}$ and $0 \leq r_{\gamma} < m_{\gamma}$. Using the uniform ergodicity in [10] and that π_{θ} is invariant for $P_{t,\theta}$ with $t \geq 0$, see Lemma 16, we obtain for all $\theta \in \Theta, \gamma \in (0, \bar{\gamma})$ and $k \in \mathbb{N}$

$$\begin{aligned} \|\pi_{\theta}R_{\gamma,\theta}^{k} - \pi_{\theta}P_{\gamma k,\theta}\|_{V^{1/2}} &\leqslant \sum_{\ell=0}^{q_{\gamma}-1} \|\pi_{\theta}P_{\gamma(\ell+1)m_{\gamma},\theta}R_{\gamma,\theta}^{(q_{\gamma}-(\ell+1))m_{\gamma}+r_{\gamma}} - \pi_{\theta}P_{\gamma\ell m_{\gamma},\theta}R_{\gamma,\theta}^{(q_{\gamma}-\ell)m_{\gamma}+r_{\gamma}}\|_{V^{1/2}} \\ &+ \|\pi_{\theta}P_{\gamma(q_{\gamma}m_{\gamma}+r_{\gamma}),\theta} - \pi_{\theta}P_{\gamma q_{\gamma}m_{\gamma},\theta}R_{\gamma,\theta}^{r_{\gamma}}\|_{V^{1/2}} \\ &\leqslant \sum_{\ell=0}^{q_{\gamma}-1} C\xi^{\gamma m_{\gamma}(q_{\gamma}-(\ell+1))}\|\pi_{\theta}P_{\gamma\ell m_{\gamma}\theta}P_{m_{\gamma}\gamma,\theta} - \pi_{\theta}P_{\gamma\ell m_{\gamma},\theta}R_{\gamma,\theta}^{m_{\gamma}}\|_{V^{1/2}} \\ &+ \|\pi_{\theta}P_{\gamma(q_{\gamma}m_{\gamma}+r_{\gamma}),\theta} - \pi_{\theta}P_{\gamma q_{\gamma}m_{\gamma},\theta}R_{\gamma,\theta}^{r_{\gamma}}\|_{V^{1/2}} \\ &\leqslant \|\pi_{\theta}P_{m_{\gamma}\gamma,\theta} - \pi_{\theta}R_{\gamma,\theta}^{m_{\gamma}}\|_{V^{1/2}} \sum_{\ell=1}^{q_{\gamma}} C\xi^{\ell\gamma m_{\gamma}} + \|\pi_{\theta}P_{\gamma r_{\gamma},\theta} - \pi_{\theta}R_{\gamma,\theta}^{r_{\gamma}}\|_{V^{1/2}}, \end{aligned}$$

where $C \ge 0, \xi \in (0, 1)$ are the constants given by [10] with drift condition $D(V^{1/2}, \lambda^{1/2}, b/2)$. Using **L1** and $(a + b)^2 \le 2(a^2 + b^2)$, we obtain that for any $x \in \mathbb{R}^d$ and $\theta \in \Theta$

$$\|\nabla_x U_{\theta}(x)\|^2 \leq 2L^2 \|x\|^2 + 2\sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\|^2 \leq 2 \left[L^2 M_{V,2} + \sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\|^2 \right] V(x) .$$
 (46)

Since \mathcal{A}_{θ} satisfies a $\mathscr{D}(V, \zeta, \beta)$ and $R_{\gamma, \theta}$ satisfies $\mathbf{D}_{d}(V, \lambda, b)$ for any $\theta \in \Theta$, we obtain that for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$

$$\pi_{\theta} P_{\gamma m_{\gamma},\theta}(V) \leqslant \tilde{A}_{1} , \qquad \pi_{\theta} R_{\gamma,\theta}^{m_{\gamma}}(V) \leqslant \tilde{A}_{1} , \qquad \tilde{A}_{1} = \beta/\zeta + b\lambda^{-\bar{\gamma}} \log(1/\lambda)^{-1} .$$
(47)

Combining (46) and (47) we can apply [10] and we get for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$

$$\|\pi_{\theta}P_{\gamma m_{\gamma},\theta} - \pi_{\theta}R_{\gamma,\theta}^{m_{\gamma}}\|_{V^{1/2}} \leqslant B_{3}^{\prime}\gamma^{1/2} , \ \|\pi_{\theta}P_{\gamma r_{\gamma},\theta} - \pi_{\theta}R_{\gamma,\theta}^{r_{\gamma}}\|_{V^{1/2}} \leqslant B_{3}^{\prime}\gamma^{1/2} , \tag{48}$$

with

$$B'_{3} = \tilde{A}_{1}^{1/2} (1+\bar{\gamma})^{1/2} \left\{ d + 2\bar{\gamma} (\mathbf{L}^{2} + \sup_{\theta \in \Theta} \|\nabla_{x} U_{\theta}(0)\|) \tilde{A}_{1} \right\}^{1/2} \mathbf{L}$$

Combining (45) and (48) we get for any $k \in \mathbb{N}$, $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$

$$\|\pi_{\theta} R_{\gamma,\theta}^{k} - \pi_{\theta} P_{\gamma k,\theta}\|_{V^{1/2}} \leqslant CB_{3}' \left(\sum_{\ell=1}^{q_{\gamma}} \xi^{\gamma m_{\gamma} \ell} + 1 \right) \gamma^{1/2} \leqslant CB_{3}' \left\{ 1 + 1/(1-\xi) \right\} \gamma^{1/2} ,$$

where we used that $\xi^{\gamma m_{\gamma}} \leq \xi$. We prove the result upon taking the limit in k.

B.2.3 Checking H2

Proposition 19. Assume **L1** and there exist a measurable function $V : \mathbb{R}^d \to (1, +\infty), \lambda \in (0, 1), b \ge 0$ and $\bar{\gamma} > 0$ such that $\lim_{\|x\|\to+\infty} V(x) = +\infty$ and for any $\theta \in \Theta$ and $\gamma \in (0, \bar{\gamma}), R_{\gamma,\theta}$ satisfies for any $x \in \mathbb{R}^d$,

$$R_{\gamma,\theta}V(x) \leqslant \lambda^{\gamma}V(x) + b\gamma$$

Then for any $p \in \mathbb{N}^*$, there exists $B_{4,p} \ge 0$ such that for any $f : \mathbb{R}^d \to \mathbb{R}$, with $\sup_{\mathbb{R}^d} \left| f/V^{1/p} \right| \le 1$, $\theta \in \Theta, \ \gamma \in (0, \overline{\gamma}), \ x \in \mathbb{R}^d$ and $n \in \mathbb{N}$

$$n^{-p} \mathbb{E}\left[\left|\sum_{k=1}^{n} f(\bar{X}_{k}) - \pi_{\theta}(f)\right|^{p}\right] \leq B_{4,p} V(x) ,$$

where $(\bar{X}_k)_{k\in\mathbb{N}}$ is the Markov chain starting from x defined by (6). In addition, $B_{4,p} = CS_b(b)S_\lambda(\lambda)$ with C a universal constant and S_b , S_λ are polynomials such that

$$\sup_{b \in \mathbb{R}} \left| (1+b)^{-(p+1)} \mathbf{S}_b(b) \right| \leq 1, \qquad \sup_{[0,1)} \left| (1-\lambda)^{p+1} \mathbf{S}_\lambda(\lambda) \right| \leq 1.$$

Proof. Let $p \in \mathbb{N}^*$ and $f : \mathbb{R}^d \to \mathbb{R}$, with $\sup_{\mathbb{R}^d} |f/V^{1/p}| \leq 1, \theta \in \Theta$ and $\gamma \in (0, \bar{\gamma})$. Using the Hölder inequality, setting q = p/(p-1), and the triangle inequality we get that

$$\left| \left(\sum_{k=1}^{n} f(\bar{X}_{k}) - \pi_{\theta}(f) \right)^{p} \right| \leq n^{p/q} \sum_{k=1}^{n} \left| f(\bar{X}_{k}) - \pi_{\theta}(f) \right|^{p} \leq n^{p/q} \sum_{k=1}^{n} \{ |f| (\bar{X}_{k}) + |\pi_{\theta}(f)| \}^{p} \leq n^{p/q} \left(\sum_{j=0}^{p} {p \choose j} \pi_{\theta}(|f|)^{p-j} \sum_{k=1}^{n} |f|^{j} (\bar{X}_{k}) \right),$$

$$(49)$$

Using Lemma 16 and $\sup_{\mathbb{R}^d} |f/V^{1/p}| \leq 1$ we obtain that

$$\pi_{\theta}(|f|) \leqslant CA_0^{1/p} , \quad \mathbb{E}_x \left[\sum_{j=1}^n |f|^j(\bar{X}_k) \right] \leqslant nCA_0^{j/p} V(x) , \quad A_0 \leqslant (1+b)(1+\lambda^{-\bar{\gamma}}/\log(1/\lambda)) \quad (50)$$

where C is a universal constant. Combining (49) and (50) we get that

$$\sup_{\|f\|_{V^{1/2}} \leqslant 1} \mathbb{E}_{x} \left[n^{-p} \left\{ \sum_{k=1}^{n} f(\bar{X}_{k}) - \pi_{\gamma,\theta}(f) \right\}^{p} \right] \leqslant B_{4,p} n^{1+p/q-p} V(x) \leqslant B_{4,p} V(x) ,$$

where $B_{4,p} = CS_b(b)S_\lambda(\lambda)$ with C a universal constant and S_b , S_λ are polynomials such that

$$\sup_{b\in\mathbb{R}} \left| (1+|b|)^{-1} \mathcal{S}_b(b) \right| \leq 1 , \qquad \sup_{\lambda\in[0,1)} \left| (1-\lambda) \mathcal{S}_\lambda(\lambda) \right| \leq 1 .$$

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B.2.4 Proof of Theorem 5

L 1 and L 2 or L 3 ensure a uniform drift condition on $R_{\gamma,\theta}$, see Proposition 12 and Proposition 13. Note that the Lyapunov functions V defined by Proposition 12 and Proposition 13 satisfy $\sup_{x \in \mathbb{R}^d} (1 + ||x||^2)/V(x) < +\infty$. Assumption H1 is then implied combining Lemma 17, Proposition 18, the result of uniform geometric ergodicity with constants independent from $\theta \in \Theta$ in [10] and that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $||H_{\theta}(x)|| \leq M_H V^{1/2}(x)$. Note that we have the following equalities in H1

$$A_1 = B_1$$
, $A_2 = B_2 M_H$, $A_3 = B_3 M_H$,

where B_2 is given in [10] and depends only on b and λ if L4 holds. Using [10] we show that B_3 depends only on b, λ, ζ, β and d and satisfies $\sup_{d \in \mathbb{R}} B_3(1+|d|)^{-1/2} < +\infty$.

We conclude that **H2** holds using Proposition 19 for p = 2 and that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $||H_{\theta}(x)|| \leq M_H V(x)^{1/2}$. We obtain that $A_4 = B_{4,2}M_H^2$.

B.3 Proof of Proposition 6

Proof. The proof is divided in two parts. First we show there exist $\eta > 0$ and $R \ge 0$ such that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, ||x|| > R,

$$U_{\theta}(x) - U_{\theta}(0) \ge \eta \|x\| .$$
(51)

Let M defined by

$$M = \sup_{\theta \in \Theta} U_{\theta}(0) + \sup_{\theta \in \Theta, \ x \in \overline{B}(0,1)} U_{\theta}(x) .$$

Note that by **L1** and since Θ is compact we have that $M < +\infty$. By contradiction, we show that there exists $R \ge 0$ such that for any $\theta \in \Theta \{x, U_{\theta}(x) \le M+1\} \subset B(0, R)$. Assume that for any $R \ge 0$, there exist $\theta_R \in \Theta$, $x_R \in \{x, U_{\theta_R}(x) \le M+1\}$ and $||x_R|| \ge R$. Then $\operatorname{Vol}(\operatorname{Conv}[\overline{B}(0, 1) \cup \{x_R\}])$ grows at least linearly in R independently from θ . On the other hand by definition of M and x_R and since **L4**(0) holds, we have for any $R \ge 0$

$$\operatorname{Vol}(\operatorname{Conv}[\overline{\mathcal{B}}(0,1)\cup\{x_R\}]) \leqslant e^{M+1} \int_{\{U_{\theta_R}(x)\leq M+1\}} e^{-U_{\theta_R}(y)} \mathrm{d}y \leqslant e^{M+1} \sup_{\theta\in\Theta} \int_{\mathbb{R}^d} e^{-U_{\theta}(y)} \mathrm{d}y \;.$$

As a consequence, we conclude that there exists $R \ge 0$, such that for any $\theta \in \Theta$, $\{x, U_{\theta}(x) \le M+1\} \subset B(0, R)$. Set $x \in \overline{B}(0, R)^c$ and y = Rx/||x||. By L4(0) we get for any $\theta \in \Theta$

$$M + 1 \leq U_{\theta}(y) \leq RU_{\theta}(x) / \|x\| + (1 - R/\|x\|)U_{\theta}(0)$$

which gives (51) with $\eta = \{M + 1 - \sup_{\theta \in \Theta} U_{\theta}(0)\}/R > 0$ by definition of M. From (51) and L1, for any $\theta \in \Theta$ there exists x_{θ}^* minimizer of U_{θ} . Thus we get for any $x \in \overline{B}(0, R)^c$

$$U_{\theta}(x) - U_{\theta}(x_{\theta}^*) \ge U_{\theta}(x) - U_{\theta}(0) \ge \eta \|x\|$$

Furthermore we have for any $x \in \mathbb{R}^d$ we have $U_{\theta}(x) - U_{\theta}(x_{\theta}^*) \ge 0$. From these results, [34, Theorem 2.1.5, Equation (2.1.7)], L1, L4(0) and (51) we obtain that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$

$$\langle \nabla_x U_\theta(x), x - x_\theta^* \rangle \ge (2\mathsf{L})^{-1} \| \nabla_x U_\theta(x) \|^2 + \eta \| x \| \mathbb{1}_{\overline{\mathrm{B}}(0,R)^c}(x) .$$

$$(52)$$

Furthermore, using $ab \leq 2^{-1}(\varepsilon a^2 + (b^2/\varepsilon))$, the Cauchy-Schwarz inequality and L4(0) we get for any $\theta \in \Theta$, $x \in \mathbb{R}^d$ and $\varepsilon > 0$

$$|\langle \nabla_x U_\theta(x), x_\theta^* \rangle| \leq \|\nabla_x U_\theta(x)\| \|x_\theta^*\| \leq \varepsilon \|\nabla_x U_\theta(x)\|^2 / 2 + {M_*}^2 / 2\varepsilon .$$
(53)

Combining (52) and (53), we obtain for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$ that

$$\begin{split} \langle \nabla_x U_\theta(x), x \rangle &= \langle \nabla_x U_\theta(x), x - x_\theta^* \rangle + \langle \nabla_x U_\theta(x), x_\theta^* \rangle \\ &\geqslant \langle \nabla_x U_\theta(x), x - x_\theta^* \rangle - |\langle \nabla_x U_\theta(x), x_\theta^* \rangle| \\ &\geqslant (1/2) (\mathsf{L}^{-1} - \varepsilon) \| \nabla_x U_\theta(x) \|^2 + \eta \| x \| \mathbb{1}_{\overline{\mathrm{B}}(0,R)^{\mathrm{c}}}(x) - M_*^2 / 2\varepsilon \;, \end{split}$$

which concludes the proof upon taking $\varepsilon = (2L)^{-1}$.

B.4 Proof of Corollary 7

Proof. Note that applying Proposition 13 and Proposition 15 with $\tilde{V}(x) = (||x|| + 1)^2/d$ we obtain that b, λ, ζ and β do not depend on the dimension d. Theorem 5 still holds for this Lyapunov function upon noticing that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $||H_{\theta}(x)|| \leq d^{1/2} M_H \tilde{V}(x)$. Let $M'_H = M_H d^{1/2}$

Theorem 3 holds using Theorem 5. We recall that the sequence $(E_n)_{n \in \mathbb{N}}$ is given by for any $n \in \mathbb{N}$ by

$$E_{n} = 2M_{\Theta}^{2} + 2M_{\Theta}\mathbb{E}\left[V(X_{0}^{0})\right] \left[\frac{(A_{1} \vee 1)A_{2}}{\xi^{\bar{\gamma}}\log(1/\xi)} \sum_{k=0}^{n} \{\delta_{k+1}/(\gamma_{k}m_{k})\} + A_{3} \sum_{k=0}^{n} \delta_{k+1}B(\gamma_{k})\right] + A_{4}\mathbb{E}\left[V(X_{0}^{0})\right] \sum_{k=0}^{n} \delta_{k+1}^{2}m_{k}^{-2}D(\gamma_{k},m_{k}).$$
(54)

Using that $A_1 = B_1$, $A_2 = B_2 M'_H$, $A_3 = B_3 M'_H$ and $A_4 = B_{4,2} M'^2_H$ we have

$$E_{n} = 2M_{\Theta}^{2} + 2d^{1/2}M_{\Theta}M_{H}\mathbb{E}\left[V(X_{0}^{0})\right] \left[\frac{(B_{1}\vee 1)B_{2}}{\xi^{\bar{\gamma}}\log(1/\xi)}\sum_{k=0}^{n}\{\delta_{k+1}/(\gamma_{k}m_{k})\} + B_{3}\sum_{k=0}^{n}\delta_{k+1}B(\gamma_{k})\right] + dB_{4,2}M_{H}^{2}\mathbb{E}\left[V(X_{0}^{0})\right]\sum_{k=0}^{n}\delta_{k+1}^{2}m_{k}^{-2}D(\gamma_{k},m_{k}).$$
(55)

Since b, λ, ζ and β do not depend on the dimension $d, B_1, B_2, B_{4,2}$ do not depend on the dimension as well. We have $B_3 \leq Cd^{1/2}$, where C is a constant independent of the dimension. Finally we obtain that E_n is bounded by Cd where C is a constant independent of the dimension.

B.5 Proof of Theorem 8

We preface the proof by a technical lemma.

Lemma 20. Let $V : \mathbb{R}^d \to [1, +\infty)$ and $M_{V,4} \ge 0$ such that $\sup_{x \in \mathbb{R}^d} (1 + ||x||^4) / V(x) \le M_{V,4}$. Let $M \ge 0$ such that for any $\theta \in \Theta$, $\gamma \in (0, \bar{\gamma})$, with $\bar{\gamma} > 0$ and $x \in \mathbb{R}^d$, $R_{\gamma,\theta}V(x) \le MV(x)$. Assume **L1** and **L5**, then we have for any $\theta_1, \theta_2 \in \Theta$, $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2$, $a \in [1/4, 1/2]$ and $x \in \mathbb{R}^d$

$$\|\delta_x R_{\gamma_1,\theta_1} - \delta_x R_{\gamma_2,\theta_2}\|_{V^a} \leqslant B_6 \left[(\gamma_2^{-1} - \gamma_1^{-1})^{1/2} + \gamma_2^{-1/2} |\gamma_1 - \gamma_2| + \gamma_2^{1/2} \|\theta_1 - \theta_2\| \right] V(x)^{2a} ,$$

where $\{R_{\gamma,\theta}, \gamma \in (0,\bar{\gamma}), \theta \in \Theta\}$ is the sequence of Markov kernels associated with the recursion (6) and

$$B_{6} = \max\left(2M^{1/2} \left[d/4 + \sup_{\theta \in \Theta} \|\nabla_{x} U_{\theta}(0)\|^{2} + L^{2} M_{4,V}^{1/2}\right]^{1/2}, (2M)^{1/2} L_{U}\right)$$

Proof. Using [12, Lemma 24] we have that for any $\theta_1, \theta_2 \in \Theta$ and $\gamma_1, \gamma_2 \in (0, \bar{\gamma})$

$$\begin{aligned} \|\delta_{x}R_{\gamma_{1},\theta_{1}} - \delta_{x}R_{\gamma_{2},\theta_{2}}\|_{V^{a}} &\leq \sqrt{2} \left(R_{\gamma_{1},\theta_{1}}V^{2a}(x) + R_{\gamma_{2},\theta_{2}}V^{2a}(x)\right)^{1/2} \operatorname{KL}\left(\delta_{x}R_{\gamma_{1},\theta_{1}}|\delta_{x}R_{\gamma_{2},\theta_{2}}\right)^{1/2} \\ &\leq 2M^{a}V^{a}(x)\operatorname{KL}\left(\delta_{x}R_{\gamma_{1},\theta_{1}}|\delta_{x}R_{\gamma_{2},\theta_{2}}\right)^{1/2} \end{aligned}$$
(56)

Denote for any $z \in \mathbb{R}^d$ and $\sigma^2 > 0$, γ_{z,σ^2} the *d*-dimensional Gaussian distribution with mean z and covariance matrix σ^2 Id. Using that for any $z_1, z_2 \in \mathbb{R}^d$ and $\sigma_1, \sigma_2 > 0$,

$$\operatorname{KL}\left(\gamma_{z,\sigma^{2}}|\gamma_{z,\sigma^{2}}\right) = \ln(\sigma_{2}^{2}/\sigma_{1}^{2}) + d\sigma_{2}^{-2}\left\{1 - \sigma_{2}^{2}/\sigma_{1}^{2}\right\}/2 + \sigma_{2}^{-2}\|\mu_{1} - \mu_{2}\|^{2}/2$$

we obtain that for any $\theta_1, \theta_2 \in \Theta, \gamma_1, \gamma_2 \in (0, \bar{\gamma})$ and $x \in \mathbb{R}^d$,

$$\operatorname{KL}\left(\delta_{x}R_{\gamma_{1},\theta_{1}}|\delta_{x}R_{\gamma_{2},\theta_{2}}\right) = \ln(\gamma_{2}/\gamma_{1}) + d\gamma_{2}^{-1}(1-\gamma_{2}/\gamma_{1})/4 + \gamma_{2}^{-1}\Xi/4 , \qquad (57)$$

where Ξ satisfies the following inequalities

$$\begin{aligned} \Xi &= \|\gamma_1 \nabla_x U_{\theta_1}(x) - \gamma_2 \nabla_x U_{\theta_2}(x)\|^2 \\ &= \|\gamma_1 \nabla_x U_{\theta_1}(x) - \gamma_2 \nabla_x U_{\theta_1}(x) + \gamma_2 \nabla_x U_{\theta_1}(x) - \gamma_2 \nabla_x U_{\theta_2}(x)\|^2 \\ &\leqslant 2 \|\gamma_1 \nabla_x U_{\theta_1}(x) - \gamma_2 \nabla_x U_{\theta_1}(x)\|^2 + 2 \|\gamma_2 \nabla_x U_{\theta_1}(x) - \gamma_2 \nabla_x U_{\theta_2}(x)\|^2 \\ &\leqslant 2 (\gamma_1 - \gamma_2)^2 \|\nabla_x U_{\theta_1}(x)\|^2 + 2\gamma_2^2 \|\nabla_x U_{\theta_1}(x) - \nabla_x U_{\theta_2}(x)\|^2 \\ &\leqslant 2 (\gamma_1 - \gamma_2)^2 \|\nabla_x U_{\theta_1}(x)\|^2 + 2\gamma_2^2 L_U^2 \|\theta_1 - \theta_2\|^2 V^{2a}(x) , \end{aligned}$$
(58)

where we used L5 in the last line. Using L5 again, $\sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\| < +\infty$ by L1, we obtain that for any $\theta \in \Theta$, $x \in \mathbb{R}^d$ and $a \in [1/4, 1/2]$ we have

$$\|\nabla_x U_{\theta}(x)\|^2 \leq 2(\|\nabla_x U_{\theta}(x) - \nabla_x U_{\theta}(0)\|^2 + \sup_{\theta \in \Theta} \|\nabla_x U_{\theta}(0)\|^2) \leq C_{\Theta} V^{2a}(x) ,$$

with $C_{\Theta} = 2 \| \sup_{\theta \in \Theta} \| \nabla_x U_{\theta}(0) \|^2 + 2 L^2 M_{4,V}^{1/2}$. Combining this result and (58) in (57), we have that for any $\theta_1, \theta_2 \in \Theta, \gamma_1, \gamma_2 \in (0, \bar{\gamma})$ with $\gamma_1 > \gamma_2$ and $x \in \mathbb{R}^d$

$$\begin{aligned} & \operatorname{KL}\left(\delta_{x}R_{\gamma_{1},\theta_{1}}|\delta_{x}R_{\gamma_{2},\theta_{2}}\right) \\ & \leq d\gamma_{2}^{-1}(1-\gamma_{2}/\gamma_{1})/4 + \gamma_{2}^{-1}(\gamma_{1}-\gamma_{2})^{2} \|\nabla_{x}U(\theta_{1},x)\|^{2}/2 + \gamma_{2}L_{U}^{2}\|\theta_{1}-\theta_{2}\|^{2}V^{2a}(x)/2 \\ & \leq \left[d\gamma_{2}^{-1}(1-\gamma_{2}/\gamma_{1})/4 + \gamma_{2}^{-1}(\gamma_{1}-\gamma_{2})^{2}C_{\Theta}/2 + \gamma_{2}L_{U}^{2}\|\theta_{1}-\theta_{2}\|^{2}/2\right]V^{2a}(x) \end{aligned}$$

This result substituted in (56) completes the proof with the inequality for any $a, b \in \mathbb{R}_+$, $(a+b)^{1/2} \leq a^{1/2} + b^{1/2}$.

Proof of Theorem 8. L1 and L2 or L3 ensure a uniform drift condition on $R_{\gamma,\theta}$, see Proposition 12 and Proposition 14. Note that the Lyapunov functions V defined by Proposition 12 and Proposition 14 satisfy $\sup_{x \in \mathbb{R}^d} (1 + ||x||^4) / V(x) < +\infty$. Assumption H3 is obtained using Lemma 17, . Combining H3, the result of uniform geometric ergodicity with constants independent from $\theta \in \Theta$ in [10] and that for any $\theta \in \Theta$ and $x \in \mathbb{R}^d$, $||H_{\theta}(x)|| \leq M_H V^{1/2}(x)$ we obtain H5. Using Proposition 18 and H5 we obtain H1-(ii). H4 is a direct consequence of Lemma 20