Optimal Scaling for the Proximal Langevin Algorithm in High Dimensions

Natesh S. Pillai *

Department of Statistics,
Harvard University
MA, USA, 02138-2901
e-mail: pillai@fas.harvard.edu

Abstract: The Metropolis-adjusted Langevin (MALA) algorithm is a sampling algorithm that incorporates the gradient of the logarithm of the target density in its proposal distribution. In an earlier joint work [PST12], the author had extended the seminal work of [RR98] and showed that in stationarity, MALA applied to an $N$-dimensional approximation of the target will take $O(N^{1/3})$ steps to explore its target measure. It was also shown in [RR98, PST12] that, as a consequence of the diffusion limit, the MALA algorithm is optimized at an average acceptance probability of 0.574. In [Per16], Pereyra introduced the proximal MALA algorithm where the gradient of the log target density is replaced by the proximal function (mainly aimed at implementing MALA non-differentiable target densities). In this paper, we show that for a wide class of twice differentiable target densities, the proximal MALA enjoys the same optimal scaling as that of MALA in high dimensions and also has an average optimal acceptance probability of 0.574. The results of this paper thus give the following practically useful guideline: for smooth target densities where it is expensive to compute the gradient while implementing MALA, users may replace the gradient with the corresponding proximal function (that can be often computed relatively cheaply via convex optimization) without losing any efficiency. We show this for two class of examples. First, for the product of Gaussians, we identify the optimal scale for proximal MALA and show that it is identical to MALA; our calculations further suggest that the same result should hold for a wide class of twice differentiable, log-concave product measures. Next, following the exact framework used in [PST12], we define a version of the proximal MALA algorithm in a Hilbert space. We show that for a certain class of twice differentiable, infinite dimensional non-product measures commonly used in applications, the proximal MALA applied to an $N$-dimensional approximation of the target also will take $O(N^{1/3})$ steps to explore the invariant measure, with an optimal acceptance probability of 0.574. This confirms some of the empirical observations made in [Per16].

AMS 2000 subject classifications: Primary 60J20; secondary 65C05, Markov Chain Monte Carlo, Metropolis Adjusted Langevin Algorithm, Scaling limit, Diffusion Approximation, Convex Optimization, Proximal Operators, Moreau Envelope.

1. Introduction

The Langevin diffusion in $\mathbb{R}^N$

$$dX_t = \nabla \log \pi^N(X_t)dt + \sqrt{2}\,dW_t$$

(1.1)

under practically realistic regularity assumptions on the measure $\pi^N$ has $\pi^N$ its invariant measure. The Langevin algorithm has been one of the workhorses for sampling probability measures; it is widely used in bayesian statistics [RC04], data assimilation, inverse problems [Stu10] and machine learning e.g., [WT11, Lam21], among other areas of data science. The time discretization of $X_t$ with step-size $\delta$ gives rise to the (unadjusted) discrete Langevin proposal:

$$y = x + \delta \nabla \log \pi^N(x) + \sqrt{2\delta}Z^N, \quad Z^N \sim N(0, I_N).$$

(1.2)

Consider a $\pi^N$-invariant Metropolis Hastings Markov chain $\{x^{k,N}\}_{k \geq 1}$ obtained by proposing $y$ from the current state $x$ according to the kernel $q(x, y)$ given by (1.2) and then accepted with probability

$$\alpha(x, y) = 1 \wedge \frac{\pi^N(y)q(y, x)}{\pi^N(x)q(x, y)}.$$ 

(1.3)
For this product measure, the seminal papers \cite{RGG97,RR98} choose the proposal variance to satisfy $\delta > \lambda > 0$ and $\ell > 1$. We choose the proposal variance $\delta$ so as to have acceptance probabilities of $O(1)$ and still having an $O(N)$ probability. The “optimal scale” for the proposal variance thus strikes a balance between making large moves and still having an $O(1)$ acceptance probability as a function of the dimension $N$. To make this heuristic precise, consider the continuous interpolant of the Markov chain $X_{t}^{k}$:

$$z(t) = \left( \frac{t}{\Delta t} - k \right) x^{k+1} + \left( k + 1 - \frac{t}{\Delta t} \right) x^{k}, \quad \text{for} \quad k\Delta t \leq t < (k+1)\Delta t.$$ \hspace{1cm} (1.5)

We choose the proposal variance to satisfy $\delta = \ell \Delta t$, with $\Delta t = N^{-\gamma}$ setting the scale in terms of dimension and the parameter $\ell$ a “tuning” parameter which is independent of the dimension $N$. We now discuss how to choose $\gamma$ and $\ell$.

Suppose that $\pi_{N}$ is the product of $N$ probability densities $\pi$,

$$\pi_{N}^{N}(x) \propto \prod_{i=1}^{N} \pi_{i}(x_{i}).$$ \hspace{1cm} (1.6)

For this product measure, the seminal papers \cite{RGG97,RR98} respectively showed that, in stationarity, the “optimal” choice for $\gamma$ that maximizes the expected jumping distance is $\gamma = 1$ for the RWM algorithm and $\gamma = \frac{1}{2}$ for the MALA. Moreover, the projection of $z^{N}$ into any single fixed coordinate direction $x_{i}$ converges weakly in $C([0,T]; \mathbb{R})$ to $z$, the scalar diffusion process

$$\frac{dz}{dt} = h(\ell)[\log \pi(z)]' + \sqrt{2h(\ell)} dW.$$ \hspace{1cm} (1.7)

Here $h(\ell) > 0$ is a constant determined by the parameter $\ell$ from the proposal variance. The quantity $h(\ell)$ has the interpretation as the “speed-measure” of the limiting diffusion \cite{RR01}. Choosing $\ell$ to maximize $h(\ell)$, thus maximizing the speed of the limiting diffusion, then yields an optimal average acceptance probability of 0.234 for the Random Walk Metropolis Algorithm and 0.574 for MALA. A remarkable feature of these results is that the optimal acceptance probabilities for these two algorithms are “universal” – they hold for a wide range of $\pi$.

The above analysis shows that the number of steps required to sample the target measure grows as $O(N)$ for RWM, but only as $O(N^{3/4})$ for MALA. This quantifies the efficiency gained by use of MALA over RWM, and in particular from employing local moves informed by the gradient of the logarithm of the target density. These theoretical analyses have inspired much further research as they give useful guidelines for implementation of MALA in high dimensions: in addition to employing an explicit scale in the proposal variance as predicted by the theory, one should “tune” the proposal variance of the RWM and MALA algorithms so as to have acceptance probabilities of 0.234 and 0.574 respectively.

### 1.1. Proximal MALA algorithm

The proximal MALA algorithm was introduced by Pereyra in \cite{Per16}. For a convex function $f : \mathbb{R}^{N} \to \mathbb{R}$, $\lambda > 0$ and $\| \cdot \|$ denoting the Euclidean norm, define the proximity operator (also called the $\lambda$-Moreau envelope, \cite{BC11}):

$$\text{Prox}_{\lambda}(x) = \arg \min_{y \in \mathbb{R}^{N}} \left( f(y) + \frac{1}{2\lambda}\|y - x\|^{2} \right).$$
The following two extreme limits are well known for proximal functions (see \cite{BC11}, Chapter 12):
\[ \lim_{\lambda \to 0} \text{Prox}^\lambda (x) = x, \quad \lim_{\lambda \to \infty} f(\text{Prox}^\lambda (x)) = \inf_{y \in \mathbb{R}^N} f(y). \]

Let \( \pi^N \) be a probability density in \( \mathbb{R}^N \) and consider it’s \( \lambda \)-Moreau approximation (see Equation (3) of \cite{Per16}):
\[ \pi^N_\lambda (x) \propto \sup_{u \in \mathbb{R}^N} \pi(u) \exp \left( - \frac{1}{2\lambda} \|u - x\|^2 \right). \]

If \( \pi^N (x) \propto \exp(-\Psi(x)) \) for a convex function \( \Psi \), we have the identity:
\[ \pi^N_\lambda (x) \propto \exp \left\{ -\Psi(\text{Prox}^\lambda_\Psi (x)) \right\} \exp \left\{ - \frac{1}{2\lambda} \|\text{Prox}^\lambda_\Psi (x) - x\|^2 \right\}. \]

In addition, if \( \Psi \) is differentiable, we also have the identity (\cite{BC11}, Equation (12.28)):
\[ \frac{1}{\lambda} (x - \text{Prox}^\lambda_\Psi (x)) = \nabla \Psi(\text{Prox}^\lambda_\Psi). \quad (1.8) \]

Equation (1.8) can be thought of as an implicit gradient. Indeed, (1.8) yields that the deterministic, implicit update equation
\[ x^{k+1,N} = x^{k,N} - \lambda \nabla \Psi(x^{k+1,N}) \]
can be written as
\[ x^{k+1,N} = \text{Prox}^\lambda_\Psi(x^{k,N}). \quad (1.9) \]

Motivated by (1.8) and (1.9), in \cite{Per16}, Pereyra introduced the following modification of the discrete Langevin proposal \(^{1} (1.2):\)
\[ y = \left( 1 - \frac{\delta}{\lambda} \right)x + \frac{\delta}{\lambda} \text{Prox}^\lambda_\Psi (x) + \sqrt{2\delta} Z^N, \quad Z^N \sim \mathcal{N}(0, I_N). \quad (1.10) \]

The proximal MALA Markov chain then proceeds via the accept-reject mechanism (1.3) using the proposal given in (1.10).

Pereyra \cite{Per16} chooses \( \delta = \lambda \) on grounds of the stability of the resulting algorithm. We also make this choice. Thus our proximal MALA proposal is given by:
\[ y = \text{Prox}^\lambda_\Psi(x) + \sqrt{2\delta} Z^N, \quad Z^N \sim \mathcal{N}(0, I_N). \quad (1.11) \]

\( \text{However, it is far from clear how the scaling should change as a joint function of the parameters } \lambda \text{ and } \delta, \text{ especially when } \Psi \text{ is not differentiable.} \)

One of the main reasons why the proximal MALA was introduced in \cite{Per16} is that the proposal (1.11) can be applied to targets even when \( \Psi \) is not differentiable: \( e.g. \), the Laplace density \( \Psi(x) = |x| \). Quoting Pereyra \cite{Per16}: “finally, similarly to other MH algorithms based on local proposals, proximal-MALA may be geometrically ergodic yet perform poorly if the proposal variance \( \delta \) is either too small or very large. Theoretical and experimental studies of MALA show that for many high-dimensional target densities the value of \( \delta \) should be set to achieve an acceptance rate of approximately 40\% – 70\% (Pillai et al. 2012).”

\(^{1}\)For notational consistency, we have set \( 2\delta = \delta' \) where \( \delta' \) is the analogous parameter in Pereyra’s definition; see Equation (9) of \cite{Per16}
1.2. Our Contributions

In this paper, we theoretically confirm the empirical observation above made by Pereyra [Per16] for a wide class of target measures. Even if the target density is differentiable, in many practical applications it may be very expensive to compute the gradient, whereas it is often cheap to compute the proximal function via convex optimization. For example in many applied models encountered in data assimilation and Bayesian inverse problems [Stu10], the target density is of the form:

$$\pi^N(\Theta | Y) \propto \exp \left( -\frac{1}{2\sigma^2} \| Y - G(\Theta) \|^2 + h(\Theta) \right)$$

where $$G : \mathbb{R}^N \to \mathbb{R}$$ is an expensive function to compute (such as the solution of a climate model obtained via solving a partial differential equation), $$\Theta$$ is a high dimensional parameter we wish to do statistical inference for, $$Y$$ is the observed data and $$\exp(h(\Theta))$$ denotes the prior distribution for $$\Theta$$. In such examples, it is even more expensive to compute the gradient of $$G$$ with respect to $$\Theta$$ and is often numerically unstable. Thus there is a natural need for developing derivative free sampling algorithms that enjoy the same optimality properties of Langevin algorithms.

In light of the scaling results for MALA mentioned before, it is thus natural to ask whether the proximal MALA algorithm enjoys similar scaling properties as that of MALA for differentiable target densities. We prove that this is indeed the case: the proximal MALA enjoys the same optimal scaling as that of MALA in high dimensions and also has an average optimal acceptance probability of 0.574. The results of our paper thus give the following practically useful guideline: for smooth target densities where the gradient is expensive to compute while implementing MALA, users may replace the gradient with the corresponding proximal function without losing any efficiency; furthermore, users can set the proximal parameter $$\delta$$ to $$N^{-1/3}$$ and tune the algorithm to have an acceptance probability of 0.574 just as in MALA. Our paper takes the first theoretical steps towards harnessing the powerful tools of convex optimization to bear fruit on optimal scaling of MCMC algorithms. We study the optimal scaling of proximal MALA in two contexts:

1. When the target measure is a product of standard Gaussians, in Theorem 2.1 we show that the optimal scale and optimal acceptance probability for the proximal MALA algorithm is identical to that of MALA. While we do not study the case when the target is a general product measure of the form (1.6), our calculations and heuristics suggest that the same result should hold for a wide class of product measures (1.6) for which $$\Psi$$ is convex and smooth.

2. For a class of infinite dimensional non-product measures studied in [PST12], we show that the optimal scaling of $$N^{-1/3}$$ for MALA as worked out in [RR98, PST12] is also optimal for the proximal MALA algorithm when the log density is convex and differentiable; see Theorem 5.1 for the formal statement of our main result.

Let us give a high-level explanation of why the proximal MALA enjoys the same scaling as that of MALA when $$\Psi$$ is differentiable. When $$\Psi$$ is smooth, it can be shown under reasonable assumptions on the second derivative of $$\Psi$$ that:

$$|\text{Prox}_\delta^\Psi(x) - x| = \mathcal{O}(\delta). \quad (1.12)$$

Consequently, setting $$\lambda = \delta$$ in the implicit Euler identity (1.8) and using (1.12) yields that

$$\text{Prox}_\delta^\Psi(x) = x - \delta \nabla \Psi(\text{Prox}_\delta^\Psi(x)) = x - \delta \nabla \Psi(x) + R(x, \delta), \quad R(x, \delta) = \mathcal{O}(\delta^2). \quad (1.13)$$

The remainder term $$R(x, \delta)$$ is $$\mathcal{O}(\delta^2)$$. Comparing this with (1.11), we see that the proximal MALA proposal can be written as

$$y = x - \delta \nabla \Psi(x) + R(x, \delta) + \sqrt{2\delta} Z^N, \quad Z^N \sim N(0, I_N) \quad (1.14)$$

where $$x_{\text{MALA}}$$ is the MALA proposal. In high dimensions, the drift term in the diffusion limit comes from $$\mathcal{O}(\delta)$$ term; the $$\mathcal{O}(\delta^2)$$ remainder term $$R(x, \delta)$$ does not contribute to the diffusion limit and vanishes in the large $$N$$ limit. Our paper formalizes this observation for a class of infinite dimensional models studied in [PST12]; refer to Equation (4.7), Lemma 7.8 and the related discussion in Section 4.1.
1.3. Infinite Dimensional Diffusions

Motivated by applications in data assimilation, inverse problems and Bayesian nonparametrics (see [Stu10] and [HSV11]), the papers [MPS12] and [PST12] extended the results of product measures [RR98] to certain infinite dimensional non-product target measures. In both of these papers, the target measure of interest, \( \pi \), is on an infinite dimensional real separable Hilbert space \( \mathcal{H} \) and is absolutely continuous with respect to a Gaussian measure \( \pi_0 \) on \( \mathcal{H} \) with mean zero and covariance operator \( C \). This framework for the analysis of MCMC in high dimensions was first studied in the papers [BS09, BRSV08, BRS09]. The Radon-Nikodym derivative defining the target measure is assumed to have the form

\[
\frac{d\pi}{d\pi_0}(x) = M_\psi \exp(-\Psi(x))
\]  

(1.15)

for a real-valued functional \( \Psi : \mathcal{H}^s \to \mathbb{R} \) defined on a subspace \( \mathcal{H}^s \subset \mathcal{H} \) that contains the support of the reference measure \( \pi_0 \); here \( M_\psi \) is a normalizing constant.

It is proved in [DPZ92, HAVW05, HSV07] that the measure \( \pi \) is invariant for \( \mathcal{H} \)-valued SDEs (or stochastic PDEs – SPDEs) with the form

\[
\frac{dz}{dt} = -h(\ell)(z + C\nabla \Psi(z)) + \sqrt{2h(\ell)} \frac{dW}{dt}, \quad z(0) = z^0
\]  

(1.16)

where \( W \) is a Brownian motion (see [DPZ92]) in \( \mathcal{H} \) with covariance operator \( C \) and any constant \( h(\ell) > 0 \).

In [PST12], the MALA algorithm was studied when applied to a sequence of finite dimensional approximations of \( \pi \) as in (1.15). The continuous time interpolant of the Markov chain \( z^N \) given by (1.5) is shown to converge weakly to \( z \) solving (1.16) in \( C([0,T]; \mathcal{H}^s) \). Furthermore, the scaling of the proposal variance which achieves this scaling limit is inversely proportional to \( N^{1/3} \) (i.e., corresponds to the exponent \( \gamma = -1/3 \)) and the speed of the limiting diffusion process is maximized at the same universal acceptance probability of 0.574 that was found for product measures [RR98].

1.4. Notation

Throughout the paper we use the following notation in order to compare sequences and to denote conditional expectations.

- Two sequences \( \{\alpha_n\} \) and \( \{\beta_n\} \) satisfy \( \alpha_n \lesssim \beta_n \) if there exists a constant \( K > 0 \) satisfying \( \alpha_n \leq K\beta_n \) for all \( n \geq 0 \). The notations \( \alpha_n \asymp \beta_n \) means that \( \alpha_n \lesssim \beta_n \) and \( \beta_n \lesssim \alpha_n \).
- Two sequences of real functions \( \{f_n\} \) and \( \{g_n\} \) defined on the same set \( D \) satisfy \( f_n \lesssim g_n \) if there exists a constant \( K > 0 \) satisfying \( f_n(x) \leq Kg_n(x) \) for all \( n \geq 0 \) and all \( x \in D \). The notations \( f_n \asymp g_n \) means that \( f_n \lesssim g_n \) and \( g_n \lesssim f_n \).
- The notation \( \mathbb{E}_x[f(x,\xi)] \) denotes expectation with respect to \( \xi \) with the variable \( x \) fixed.

2. A Simple Example: Product of Gaussians

We start with a simple case, where the target measure is the product of standard Gaussians:

\[
\pi^N(x) \propto \prod_{i=1}^{N} \exp(-x_i^2/2).
\]  

(2.1)

The MALA proposal for \( \pi^N \) given in (2.1) is:

\[
y = x(1-\delta) + \sqrt{2\delta} Z, \quad Z \sim N(0,I_N).
\]

The Metropolis-Hastings acceptance ratio \( \alpha(x,y) \) given in (1.3) with

\[
q(x,y) = \prod_{i=1}^{N} \exp \left( -\frac{1}{4\delta} (y_i - x_i(1-\delta))^2 \right).
\]
Thus if we set $\delta q$ with the corresponding $u$ and the maximizer $\hat{a}$
Moreover, the variance of the $O(\delta^{3/2})$ yields:
Since the chain is at stationarity, the first three summands in (2.2) have expectations zero:
Moreover, the variance of the $O(\delta^{3/2})$ satisfies:
Thus if we set $\delta = \ell N^{-1/3}$, using the fact that $\frac{1}{N} \sum_{i=1}^{N} x_i^2 \to 1$ almost surely, we obtain that
and the acceptance probability:
In particular, $L_n = O(1)$ for $\delta = N^{-1/3}$, and thus the optimal scale that makes the size of acceptance probability equal to $O(1)$ corresponds to $\delta = N^{-1/3}$. The ongoing computation generalizes for quite a large class of product measures $\pi^N$ far beyond Gaussians, and forms the basis of the diffusion limit obtained in [RR98]. Finally, to have the optimal acceptance probability of $0.574$ that maximizes the speed of the limiting diffusion, all one needs is that the limiting Gaussian random variable $Z_\ell$ satisfy:
Indeed, once we have the relation (2.4), the limiting diffusion has the speed measure:
for some constant $K$ that depends on the target measure and $\Phi$ is the CDF of the standard Gaussian distribution. As shown in Theorem 2 of [RR98], the value of $\ell$ that maximizes $h(\ell)$ is independent of $K$ since making the transformation $u = \frac{K}{2} \ell^3$ yields that
and the maximizer $\hat{u}$ of the latter term is independent of $K$, see Theorem 2 of [RR98]. Thus the optimal acceptance probability is also independent of $K$: it is just $a = 2\Phi(-\hat{u})$.
Next, we perform the same computation for the proximal MALA algorithm. The proximal MALA proposal for $\pi^N$ given in (2.1) is:
with the corresponding $q(x, y)$:

\begin{equation}
q(x, y) = \prod_{i=1}^{N} \exp\left( -\frac{1}{4\delta} (y_i - \frac{x_i}{(1 + \delta)})^2 \right).
\end{equation}

\footnote{We used MATHEMATICA for obtaining this expansion; also see [RR98].}
Theorem 2.1. For the proximal MALA proposal given in (2.5), the choice of \( \delta = \ell N^{-1/3} \) yields an acceptance probability of \( O(1) \). The limiting acceptance probability \( a(\ell) \) can be expressed as \( a(\ell) = \mathbb{E}(1 \wedge e^{\delta t}) \) where \( \tilde{Z}_t \) is a Gaussian variable satisfying (2.4).

Proof. As before, expanding \( L_n = \log \left( \frac{\pi^N(y|x,y)}{\pi^N(x|y,x)} \right) \) in terms of \( \delta \) yields:

\[
L_n = -\frac{3}{\sqrt{2}} \delta^{3/2} \sum_{i=1}^{N} x_i Z_i + \frac{3}{2} \delta^2 \sum_{i=1}^{N} (x_i^2 - Z_i^2) + \delta^{5/2} \sum_{i=1}^{N} x_i Z_i + \frac{1}{4} \delta^3 \sum_{i=1}^{N} (8z_i^2 - 17x_i^2) + O \left( \delta^{7/2} \right). \tag{2.6}
\]

Again, using the fact that the chain is at stationarity, we see that the summands of \( \delta^{3/2}, \delta^2 \) and \( \delta^{5/2} \) in the expansion (2.6) all have mean zero. Furthermore, for the choice of \( \delta = \ell N^{-1/3} \), we have \( L_n \Rightarrow \tilde{Z}_t \) with \(-\frac{\delta^2}{2} = 2\mathbb{E}(\tilde{Z}_t) = \text{Var}(\tilde{Z}_t) \) satisfying (2.4), and the proof is finished. \( \square \)

While we do not prove a diffusion limit, the arguments laid out in Section 1.2 can be used to prove a diffusion limit for any single component of the piecewise interpolant of the proximal Markov chain described above. Consequently, Theorem 2.1 yields that the optimal acceptance probability for proximal MALA algorithm is also 0.574 in the case where the target measure is the product of Gaussians.

Remark 2.2. While Theorem 2.1 is only worked out for product of Gaussians, the result strongly suggests that the same optimal scale and acceptance probability should hold for a large class of measures obtained as products of log-concave target densities. The main reason for this is that the optimal scale and optimal acceptance probability results are “universal”; the specifics of target distributions should not matter. In particular, the Gaussian distribution (as used in Theorem 2.1) plays no special role in optimality of MALA and nor should play a role here. We focused on this case for clarity of exposition.

3. Infinite Dimensional Target Measure

We keep the framework in this paper very close to that of [PST12] so that the reader can easily compare our results to that of the MALA algorithm obtained in that paper. Let \( \mathcal{H} \) be a separable Hilbert space of real valued functions with scalar product denoted by \( \langle \cdot, \cdot \rangle \) and associated norm \( \| x \| = \langle x, x \rangle \). Consider a Gaussian probability measure \( \pi_0 \) on \((\mathcal{H}, \| \cdot \|)\) with covariance operator \( \mathcal{C} \). The general theory of Gaussian measures [DPZ92] ensures that the operator \( \mathcal{C} \) is positive and trace class. Let \( \{ \varphi_j, \lambda_j^2 \}_{j \geq 1} \) be the eigenfunctions and eigenvalues of the covariance operator \( \mathcal{C} \):

\[ \mathcal{C} \varphi_j = \lambda_j^2 \varphi_j, \quad j \geq 1. \]

We assume a normalization under which the family \( \{ \varphi_j \}_{j \geq 1} \) forms a complete orthonormal basis in the Hilbert space \( \mathcal{H} \), which we refer to us as the Karhunen-Loève basis. Any function \( x \in \mathcal{H} \) can be represented in this basis via the expansion

\[ x = \sum_{j=1}^{\infty} x_j \varphi_j, \quad x_j \overset{\text{def}}{=} \langle x, \varphi_j \rangle. \tag{3.1} \]

Throughout this paper we will often identify the function \( x \) with its coordinates \( \{ x_j \}_{j=1}^{\infty} \in \ell^2 \) in this eigenbasis, moving freely between the two representations. The Karhunen-Loève expansion (see [DPZ92], section White Noise expansions), refers to the fact that a realization \( x \) from the Gaussian measure \( \pi_0 \) can be expressed by allowing the coordinates \( \{ x_j \}_{j \geq 1} \) in (3.1) to be independent random variables distributed as \( x_j \sim N(0, \lambda_j^2) \). Thus, in the coordinates \( \{ x_j \}_{j \geq 1} \), the Gaussian reference measure \( \pi_0 \) has a product structure.
For every \( x \in \mathcal{H} \) we have the representation (3.1). Using this expansion, we define Sobolev-like spaces \( \mathcal{H}^r, r \in \mathbb{R} \), with the inner-products and norms defined by

\[
\langle x, y \rangle_r \overset{\text{def}}{=} \sum_{j=1}^{\infty} j^{2r} x_j y_j, \quad \| x \|^2_r \overset{\text{def}}{=} \sum_{j=1}^{\infty} j^{2r} x_j^2.
\]

Notice that \( \mathcal{H}^0 = \mathcal{H} \) and \( \mathcal{H}^r \subset \mathcal{H} \subset \mathcal{H}^{-r} \) for any \( r > 0 \). The Hilbert-Schmidt norm \( \| \cdot \|_C \) associated to the covariance operator \( C \) is defined as

\[
\| x \|^2_C = \sum_j \lambda_j^{-2} x_j^2.
\]

For \( x, y \in \mathcal{H}^r \), the outer product operator in \( \mathcal{H}^r \) is the operator \( x \otimes_{\mathcal{H}^r} y : \mathcal{H}^r \to \mathcal{H}^r \) defined by \( (x \otimes_{\mathcal{H}^r} y) z = \langle x, z \rangle_r y \) for every \( z \in \mathcal{H}^r \). For \( r \in \mathbb{R} \), let \( B_r : \mathcal{H} \to \mathcal{H} \) denote the operator which is diagonal in the basis \( \{ \varphi_j \}_{j \geq 1} \) with diagonal entries \( j^{2r} \). The operator \( B_r \) satisfies \( B_r \varphi_j = j^{2r} \varphi_j \) so that \( B_r^{1/2} \varphi_j = j^r \varphi_j \). The operator \( B_r \) lets us alternate between the Hilbert space \( \mathcal{H} \) and the Sobolev spaces \( \mathcal{H}^r \) via the identities \( \langle x, y \rangle_r = \langle B_r^{1/2} x, B_r^{1/2} y \rangle \). Since \( \| B_r^{-1/2} \varphi_k \|_r = \| \varphi_k \| = 1 \), we deduce that \( \{ B_r^{-1/2} \varphi_k \}_{k \geq 0} \) forms an orthonormal basis for \( \mathcal{H}^r \). For a positive, self-adjoint operator \( D : \mathcal{H} \to \mathcal{H} \), we define its trace in \( \mathcal{H}^r \) by

\[
\text{Tr}_{\mathcal{H}^r}(D) \overset{\text{def}}{=} \sum_{j=1}^{\infty} \langle (B_r^{-1/2} \varphi_j), D(B_r^{-1/2} \varphi_j) \rangle_r.
\]

Since \( \text{Tr}_{\mathcal{H}^r}(D) \) does not depend on the orthonormal basis, the operator \( D \) is said to be trace class in \( \mathcal{H}^r \) if \( \text{Tr}_{\mathcal{H}^r}(D) < \infty \) for some, and hence any, orthonormal basis of \( \mathcal{H}^r \). Let us define the operator \( C_r \overset{\text{def}}{=} B_r^{1/2} \mathcal{C} B_r^{1/2} \). Notice that \( \text{Tr}_{\mathcal{H}^r}(C_r) = \sum_{j=1}^{\infty} \lambda_j^2 j^{2r} \). In [PST12] it is shown that under the condition

\[
\text{Tr}_{\mathcal{H}^r}(C_r) < \infty,
\]

the support of \( \pi_0 \) is included in \( \mathcal{H}^r \) in the sense that \( \pi_0 \)-almost every function \( x \in \mathcal{H} \) belongs to \( \mathcal{H}^r \). Furthermore, the induced distribution of \( \pi_0 \) on \( \mathcal{H}^r \) is identical to that of a centered Gaussian measure on \( \mathcal{H}^r \) with covariance operator \( C_r \). For example, if \( \xi \overset{\text{d}}{\sim} \pi_0 \), then \( \mathbb{E}[\langle \xi, u \rangle_r, \langle \xi, v \rangle_r] = \langle u, C_r v \rangle_r \) for any functions \( u, v \in \mathcal{H}^r \). Thus in what follows, we alternate between the Gaussian measures \( N(0, C) \) on \( \mathcal{H} \) and \( N(0, C_r) \) on \( \mathcal{H}^r \), for those \( r \) for which (3.4) holds.

3.1. Change of Measure

Our goal is to sample from a measure \( \pi \) defined through the change of probability formula (1.15). As described above, the condition \( \text{Tr}_{\mathcal{H}^r}(C_r) < \infty \) implies that the measure \( \pi_0 \) has full support on \( \mathcal{H}^r \), i.e., \( \pi_0(\mathcal{H}^r) = 1 \). Consequently, if \( \text{Tr}_{\mathcal{H}^r}(C_r) < \infty \), the functional \( \Psi(\cdot) \) needs only to be defined on \( \mathcal{H}^r \) in order for the change of probability formula (1.15) to be valid. In this section we give assumptions on the decay of the eigenvalues of the covariance operator \( C \) of \( \pi_0 \) that ensure the existence of a real number \( s > 0 \) such that \( \pi_0 \) has full support on \( \mathcal{H}^s \). The functional \( \Psi(\cdot) \) is assumed to be defined on \( \mathcal{H}^s \) and we impose regularity assumptions on \( \Psi(\cdot) \) that ensure that the probability distribution \( \pi_0 \) is not too different from \( \pi_0 \), when projected into directions associated with \( \varphi_j \) for \( j \) large. For each \( x \in \mathcal{H}^s \) the derivative \( \nabla \Psi(x) \) is an element of the dual \( (\mathcal{H}^s)^* \) of \( \mathcal{H}^s \) comprising linear functionals on \( \mathcal{H}^s \). However, we may identify \( (\mathcal{H}^s)^* \) with \( \mathcal{H}^{-s} \) and view \( \nabla \Psi(x) \) as an element of \( \mathcal{H}^{-s} \) for each \( x \in \mathcal{H}^s \). With this identification, the following identity holds

\[
\| \nabla \Psi(x) \|_{\mathcal{L}(\mathcal{H}^r, \mathcal{R})} = \| \nabla \Psi(x) \|_{\mathcal{L}(\mathcal{H}^r, \mathcal{H}^{-s})}
\]

and the second derivative \( \partial^2 \Psi(x) \) can be identified as an element of \( \mathcal{L}(\mathcal{H}^r, \mathcal{H}^{-s}) \). To avoid technicalities we assume that \( \Psi(\cdot) \) is quadratically bounded, with first derivative linearly bounded and second derivative globally bounded. Weaker assumptions could be dealt with by use of stopping time arguments.
Assumptions 3.1. The covariance operator \( C \) and functional \( \Psi \) satisfy the following:

1. **Decay of Eigenvalues** \( \lambda_j^2 \) of \( C \): there is an exponent \( \kappa > \frac{1}{2} \) such that

\[
\lambda_j \asymp j^{-\kappa}. \tag{3.5}
\]

2. **Assumptions on \( \Psi \):** The function \( \Psi \) is convex. There exist constants \( M_i \in \mathbb{R}, i \leq 4 \) and \( s \in [0, \kappa - 1/2] \) such that for all \( x \in \mathcal{H}^s \) the functional \( \Psi : \mathcal{H}^r \rightarrow \mathbb{R} \) satisfies

\[
M_1 \leq \Psi(x) \leq M_2 \left( 1 + \|x\|_s^2 \right) \tag{3.6}
\]

\[
\|\nabla \Psi(x)\|_{-s} \leq M_3 \left( 1 + \|x\|_s \right) \tag{3.7}
\]

\[
\|\partial^2 \Psi(x)\|_{L(\mathcal{H}^s,\mathcal{H}^{-r})} \leq M_4. \tag{3.8}
\]

**Remark 3.2.** The convexity of \( \Psi \) is not assumed in [PST12]. It is not required for the MALA algorithm. In this paper we assume the convexity of \( \Psi \) so as to get a unique value for the proximal operator. This assumption is not strictly necessary for our methods to go through. However, since our key aim is to formalize the observation made in (1.14), we avoid additional complications.

**Remark 3.3.** The condition \( \kappa > \frac{1}{2} \) ensures that the covariance operator \( C \) is trace class in \( \mathcal{H} \). In fact, Equation (3.4) shows that \( C_r \) is trace-class in \( \mathcal{H}^r \) for any \( r < \kappa - \frac{1}{2} \). It follows that \( \pi_0 \) has full measure in \( \mathcal{H}^r \) for any \( r \in [0, \kappa - 1/2) \). In particular \( \pi_0 \) has full support on \( \mathcal{H}^s \).

**Remark 3.4.** The functional \( \Psi(x) = \frac{1}{2}\|x\|_s^2 \) satisfies Assumptions 3.1. It is convex, defined on \( \mathcal{H}^s \) and its derivative at \( x \in \mathcal{H}^s \) is given by \( \nabla \Psi(x) = \sum_{j \geq 0} j^{2s} x_j \varphi_j \in \mathcal{H}^{-s} \) with \( \|\nabla \Psi(x)\|_{-s} = \|x\|_s \). The second derivative \( \partial^2 \Psi(x) \in L(\mathcal{H}^s, \mathcal{H}^{-r}) \) is the linear operator that maps \( u \in \mathcal{H}^s \) to \( \sum_{j \geq 0} j^{2s} \langle u, \varphi_j \rangle \varphi_j \in \mathcal{H}^s \): its norm satisfies \( \|\partial^2 \Psi(x)\|_{L(\mathcal{H}^s, \mathcal{H}^{-r})} = 1 \) for any \( x \in \mathcal{H}^s \).

### 3.2. Finite Dimensional Approximation

We are interested in finite dimensional approximations of the probability distribution \( \pi \). To this end, we introduce the vector space spanned by the first \( N \) eigenfunctions of the covariance operator,

\[
X^N \overset{def}{=} \text{span}\{ \varphi_1, \varphi_2, \ldots, \varphi_N \}.
\]

Notice that \( X^N \subset \mathcal{H}^r \) for any \( r \in [0, +\infty) \). In particular, \( X^N \) is a subspace of \( \mathcal{H}^s \). Next, we define \( N \)-dimensional approximations of the functional \( \Psi(\cdot) \) and of the reference measure \( \pi_0 \). To this end, we introduce the orthogonal projection on \( X^N \) denoted by \( P^N : \mathcal{H}^r \rightarrow X^N \subset \mathcal{H}^s \). The functional \( \Psi(\cdot) \) is approximated by the functional \( \Psi^N : X^N \rightarrow \mathbb{R} \) defined by

\[
\Psi^N \overset{def}{=} \Psi \circ P^N. \tag{3.9}
\]

The approximation \( \pi_0^N \) of the reference measure \( \pi_0 \) is the Gaussian measure on \( X^N \) given by the law of the random variable

\[
\pi_0^N \overset{D}{=} \sum_{j=1}^N \lambda_j \xi_j \varphi_j = (\mathcal{C}^N)^{\frac{1}{2}} \xi^N
\]

where \( \xi_j \) are i.i.d standard Gaussian random variables, \( \xi^N = \sum_{j=1}^N \xi_j \varphi_j \) and \( \mathcal{C}^N = P^N \circ C \circ P^N \). Consequently we have \( \pi_0^N = N(0, \mathcal{C}^N) \). Finally, one can define the approximation \( \pi^N \) of \( \pi \) by the change of probability formula

\[
\frac{d\pi^N}{d\pi_0}(x) = M_{\Psi^N} \exp \left\{ -\Psi^N(x) \right\} \tag{3.10}
\]
where $M_{\Psi^N}$ is a normalization constant. Notice that the probability distribution $\pi^N$ is supported on $X^N$ and has Lebesgue density $\pi^N$ on $X^N$ equal to

$$\pi^N(x) \propto \exp \left(-\frac{1}{2}\|x\|^2_{C^N} - \Psi^N(x)\right). \tag{3.11}$$

In formula (3.11), the Hilbert-Schmidt norm $\| \cdot \|_{CN}$ on $X^N$ is given by the scalar product $\langle u, v \rangle_{CN} = \langle u, (C^N)^{-1}v \rangle$ for all $u, v \in X^N$. The operator $C^N$ is invertible on $X^N$ because the eigenvalues of $C$ are assumed to be strictly positive. The quantity $C^N\nabla \log \pi^N(x)$ is repeatedly used in the text and in particular appears in the function $\mu^N(x)$ given by

$$\mu^N(x) = -\left(P^N x + C^N \nabla \Psi^N(x)\right) \tag{3.12}$$

which, up to an additive constant, is $C^N\nabla \log \pi^N(x)$. This function is the drift of an ergodic Langevin diffusion that leaves $\pi^N$ invariants. Similarly, one defines the function $\mu : \mathcal{H}^s \to \mathcal{H}^s$ given by

$$\mu(x) = -\left(x + C^N \nabla \Psi(x)\right) \tag{3.13}$$

which can informally be seen as $C^N \nabla \log \pi(x)$, up to an additive constant. In Lemmas 4.1 and 4.3 of [PST12], it is shown that for $\pi_0$-almost every function $x \in \mathcal{H}$, we have $\lim_{N \to \infty} \mu^N(x) = \mu(x)$; see Section 7.1 below. This quantifies the manner in which $\mu^N(\cdot)$ is an approximation of $\mu(\cdot)$.

The next lemma gathers various regularity estimates on the functional $\Psi(\cdot)$ and $\Psi^N(\cdot)$ that are repeatedly used in the sequel. These are simple consequences of Assumptions 3.1 and proofs can be found in [MPS12] and [PST12].

**Lemma 3.5. (Properties of $\Psi$)** Let the functional $\Psi(\cdot)$ satisfy Assumptions 3.1 and consider the functional $\Psi^N(\cdot)$ defined by Equation (3.9). The following estimates hold.

1. The functionals $\Psi^N : \mathcal{H}^s \to \mathbb{R}$ satisfy the same conditions imposed on $\Psi$ given by Equations (3.6), (3.7) and (3.8) with constants that can be chosen independent of $N$.
2. The function $C^N \nabla \Psi : \mathcal{H}^s \to \mathcal{H}^s$ is globally Lipschitz on $\mathcal{H}^s$; there exists a constant $M_5 > 0$ such that

$$\|C^N \nabla \Psi(x) - C^N \nabla \Psi(y)\|_s \leq M_5 \|x - y\|_s \quad \forall x, y \in \mathcal{H}^s.$$ 

Moreover, the functions $C^N \nabla \Psi^N : \mathcal{H}^s \to \mathcal{H}^s$ also satisfy this estimate with a constant that can be chosen independent of $N$.
3. The functional $\Psi(\cdot) : \mathcal{H}^s \to \mathbb{R}$ satisfies a second order Taylor formula. There exists a constant $M_6 > 0$ such that

$$\Psi(y) - (\Psi(x) + \langle \nabla \Psi(x), y - x \rangle) \leq M_6 \|x - y\|_s^2 \quad \forall x, y \in \mathcal{H}^s. \tag{3.14}$$

Moreover, the functionals $\Psi^N(\cdot)$ also satisfy the above estimates with a constant that can be chosen independent of $N$.

**Remark 3.6.** The regularity Lemma 3.5 shows in particular that the function $\mu : \mathcal{H}^s \to \mathcal{H}^s$ defined by (3.13) is globally Lipschitz on $\mathcal{H}^s$. Similarly, it follows that $C^N \nabla \Psi^N : \mathcal{H}^s \to \mathcal{H}^s$ and $\mu^N : \mathcal{H}^s \to \mathcal{H}^s$ given by (3.12) are globally Lipschitz with Lipschitz constants that can be chosen uniformly in $N$.

4. The proximal MALA in Hilbert space

In this section, we construct a version of the proximal MALA algorithm of Pereyra [Per16] in the Hilbert space $\mathcal{H}^s$. The proximal operators are well defined in an infinite dimensional Hilbert space. The reader is

---

3For ease of notation we do not distinguish between a measure and its density, nor do we distinguish between the representation of the measure in $X^N$ or in coordinates in $\mathbb{R}^N$.

4We extend $\langle \cdot, \cdot \rangle$ from an inner-product on $H$ to the dual pairing between $H^{-s}$ and $H^s$. 

---
referred to [BC11] for a book length treatment. For the function \( \Psi : \mathcal{H}^s \mapsto (-\infty, \infty] \) and \( \lambda > 0 \), define the proximal function

\[
\text{Prox}_\lambda^\Psi(x) = \arg \min_{y \in \mathcal{H}^s} \left( \Psi(y) + \frac{1}{2\lambda} \|x - y\|_s^2 \right).
\] (4.1)

Since \( \Psi \) is convex, Proposition 12.15 of [BC11] yields that \( \text{Prox}_\lambda^\Psi(x) \) is convex and differentiable. Moreover the minimizer in (4.1) is unique due to the convexity of \( \Psi \). We also have the identity ([BC11], Corollary 17.6):

\[
\frac{1}{\lambda}(x - \text{Prox}_\lambda^\Psi(x)) = \nabla \Psi(\text{Prox}_\lambda^\Psi(x)).
\] (4.2)

The proximal functions \( \text{Prox}_\lambda^\Psi \) are defined analogously.

### 4.1. The Proximal-MALA Algorithm

Recall from (3.11) that our target measure is

\[
\pi^N(x) \propto \exp \left( -\frac{1}{2} \|x\|_C^2 - \Psi^N(x) \right).
\]

Our algorithm is motivated by the fact that the probability measure \( \pi^N \) defined by Equation (3.10) is invariant with respect to the Langevin diffusion process

\[
\frac{dz}{dt} = C^N \nabla \log \pi^N(z) + \sqrt{2} dW^N dt
\]

where \( W^N \) is a Brownian motion in \( \mathcal{H}^s \) with covariance operator \( C^N \) and \( \mu^N \) is as defined in (3.12). A natural analogue to Pereyra’s proximal proposal given in (1.11) is:

\[
y = x_{\text{Prox-MALA}} = x + (1 - \delta - C^N) x + C^N \text{Prox}_\delta^\Psi^N(x) + \sqrt{2\delta} (C^N)^{\frac{1}{2}} \xi^N \quad \text{where} \quad \delta = \ell N^{-\frac{1}{2}}. \quad (4.4)
\]

The intuition behind the proposal defined in (4.4) is the following. Applying (4.2) with \( \Psi = \Psi^N \) and \( \lambda = \delta \), we obtain that

\[
\text{Prox}_\delta^\Psi^N(x) = x - \delta \nabla \Psi^N(\text{Prox}_\delta^\Psi^N(x)) \\
\approx x - \delta \nabla \Psi^N(x) + O(\delta^2).
\]

Consequently, on \( X^N \),

\[
(1 - \delta - C^N) x + C^N \text{Prox}_\delta^\Psi^N(x) \approx x - \delta(\text{Prox}_\delta^\Psi^N(x)) \\
= x + \delta \mu^N(x). \quad (4.5)
\]

Let

\[
x_{\text{MALA}} = x + \delta \mu^N(x) + \sqrt{2\delta} (C^N)^{\frac{1}{2}} \xi^N \quad \text{where} \quad \delta = \ell N^{-\frac{1}{2}} \quad (4.6)
\]

denote the usual MALA proposal obtained from the Euler discretization of the infinite dimensional diffusion (4.3). Notice that \( (C^N)^{\frac{1}{2}} \xi^N \overset{D}{\sim} \mathcal{N}(0, C^N) \). The calculation done in (4.5) shows that our proximal MALA proposal (4.4) closely tags the MALA proposal:

\[
x_{\text{Prox-MALA}} = x_{\text{MALA}} + R^N(x, \delta) \quad (4.7)
\]
where the term
\[ R^N(x, \delta) \equiv \delta C^N \left( \text{Prox}_\delta^N (x) - x \right) \] (4.8)
can be thought of as the added “error” induced by the proximal MALA proposal as compared to MALA. As shown in Lemma 7.2, we have \( |R^N(x, \delta)| \leq \delta^2 (1 + \|x\|_s) = O(\delta^2) \). As in the product measure case, for optimal scaling only terms of \( O(\delta^{3/2}) \) and lower order contribute; thus the contribution from this remainder term to the scaling drops out in the large \( N \) limit. Consequently, the optimal scaling and the diffusion limits for the proximal MALA algorithm follows from the corresponding results for the MALA algorithm.

For streamlining further calculations, we will write the \( \mathcal{x}_{\text{Prox-MALA}} \) proposal from (4.4) as
\[ y = x + \delta \mu^N(x) + R^N(x, \delta) + \sqrt{2\delta} \left( C^N \right)^{1/2} \xi^N \quad \text{where} \quad \delta = \ell N^{-\frac{1}{4}}. \] (4.9)

### 4.2. Time evolution of the proximal MALA chain

We introduce a related parameter
\[ \Delta t := \ell^{-1} \delta = N^{-\frac{1}{4}} \]
which will be the natural time-step for the limiting diffusion process derived from the proposal above, after inclusion of an accept-reject mechanism. The scaling of \( \Delta t \), and hence \( \delta \), with \( N \) will ensure that the average acceptance probability is \( O(1) \) as \( N \) grows.

Following [PST12], we will study the Markov chain \( x^N = \{x^{k,N}\}_{k \geq 0} \) resulting from Metropolizing the proximal proposal (4.9) when it is started at stationarity: the initial position \( x^{0,N} \) is distributed as \( \pi^N \) and thus lies in \( X^N \). Therefore, the Markov chain evolves in \( X^N \); as a consequence, only the first \( N \) components of an expansion in the eigenbasis of \( \mathcal{C} \) are nonzero and the algorithm can be implemented in \( \mathbb{R}^N \). However the analysis is cleaner when written in \( X^N \subset \mathcal{H}^s \). The acceptance probability only depends on the first \( N \) coordinates of \( x \) and \( y \) and has the form
\[ \alpha^N(x, \xi^N) = 1 \wedge \pi^N(\gamma^N(y) T^N(y, x)) / \pi^N(x, y) = 1 \wedge e^{Q^N(x, \xi^N)} \] (4.10)
where the proposal \( y \) is given by Equation (4.9). The function \( T^N(\cdot, \cdot) \) is the density of the Langevin proposals (4.9) and is given by
\[ T^N(x, y) \propto \exp \left\{ -\frac{1}{2\delta} \| y - x - \delta \mu^N(x) - R^N(x, \delta) \|_{C^N}^2 \right\}. \]
The local mean acceptance probability \( \alpha^N(x) \) is defined by
\[ \alpha^N(x) = \mathbb{E}_x [\alpha^N(x, \xi^N)]. \] (4.11)
It is the expected acceptance probability when the algorithm stands at \( x \in \mathcal{H} \). The Markov chain \( x^N = \{x^{k,N}\}_{k \geq 0} \) can also be expressed as
\[ \begin{align*}
    y^{k,N} & = x^{k,N} + \delta \mu^N(x^{k,N}) + R^N(x^{k,N}, \delta) + \sqrt{2\delta} \left( C^N \right)^{1/2} \xi^{k,N} \\
    x^{k+1,N} & = \gamma^{k,N} y^{k,N} + (1 - \gamma^{k,N}) x^{k,N}
\end{align*} \] (4.12)
where \( \xi^{k,N} \) are i.i.d samples distributed as \( \xi^N \) and \( \gamma^{k,N} = \gamma^N(x^{k,N}, \xi^{k,N}) \) creates a Bernoulli random sequence with \( k^{th} \) success probability \( \alpha^N(x^{k,N}, \xi^{k,N}) \). We may view the Bernoulli random variable as
\[ \gamma^{k,N} = \mathbb{1}_{\{U^k \leq \alpha^N(x^{k,N}, \xi^{k,N})\}} \]
where \( U^k \sim \text{Uniform}(0, 1) \) is independent from \( x^{k,N} \) and \( \xi^{k,N} \).

In summary, the Markov chain that we have described in \( \mathcal{H}^s \) is, when projected onto \( X^N \), equivalent to a proximal MALA algorithm on \( \mathbb{R}^N \) for the Lebesgue density (3.11). Recall that the target measure \( \pi \) in (1.15) is the invariant measure of the SPDE (1.16). Our goal is to obtain an invariance principle for the continuous interpolant (1.5) of the Markov chain \( x^N = \{x^{k,N}\}_{k \geq 0} \) started in stationarity, \( i.e. \), to show weak convergence in \( C([0, T]; \mathcal{H}^s) \) of \( z^N(t) \) to the solution \( z(t) \) of the SPDE (1.16), as the dimension \( N \rightarrow \infty \).
5. Main Result

In this section, we present the main result of this paper. Consider the constant $\alpha(\ell) = \mathbb{E}[1 \wedge e^{Z_{\ell}}]$ where $Z_{\ell} \overset{\mathcal{D}}{\sim} \mathcal{N}(-\frac{\ell^3}{3}, \frac{\ell^3}{2})$ and define the speed function

$$h(\ell) = \ell \alpha(\ell). \quad (5.1)$$

The quantity $\alpha(\ell)$ represents the limiting expected acceptance probability of the MALA algorithm while $h(\ell)$ is the asymptotic speed function of the limiting diffusion.

**Theorem 5.1.** Let the initial condition $x^{0,N}$ of the proximal MALA algorithm be such that $x^{0,N} \sim \pi^N$ and let $z^N(t)$ be a piecewise linear, continuous interpolant of the proximal MALA algorithm (4.12) with $\Delta t = N^{-1/3}$. Then, for any $T > 0$, $z_N(t)$ converges weakly in $C([0,T], \mathcal{H}^s)$ to the diffusion process $z(t)$ given by

$$\frac{dz}{dt} = -h(\ell)(z + C\nabla \Psi(z)) + \sqrt{2h(\ell)} \frac{dW}{dt}, \quad z(0) = z^0 \sim \pi \quad (5.2)$$

with the constant $h(\ell)$ as given in (5.1). Choosing $\ell$ so as to maximize the speed function $h(\ell)$ leads to the acceptance probability of 0.574 for the proximal MALA algorithm.

**Remark 5.2.** The fact that choosing $\ell$ so as to maximize the speed function $h(\ell)$ leads to the optimal universal acceptance probability of 0.574 is known since [RR98], and is also shown in [PST12]. Thus to prove Theorem 5.1, we need only establish the diffusion limit.

5.1. Proof Strategy

The acceptance probability of the proposal (4.9) is equal to $\alpha^N(x, \xi^N) = 1 \wedge e^{Q^N(x,\xi^N)}$ and the quantity $\alpha^N(x) = \mathbb{E}_x[\alpha^N(x,\xi^N)]$ given by (4.11) represents the mean acceptance probability when the Markov chain $x^N$ stands at $x$. Recall the quantity $Q^N$ in Equation (4.10). This quantity may be expressed as

$$Q^N(x, \xi^N) = -\frac{1}{2}\left\{\|y\|^2_{\xi^N} - \|x\|^2_{\xi^N} - \left(\Psi^N(y) - \Psi^N(x)\right)\right\}$$

$$-\frac{1}{4\delta}\left\{\|x - y - \delta\mu^N(y) - R^N(y, \delta)\|^2_{\xi^N} - \|y - x - \delta\mu^N(x) - R^N(x, \delta)\|^2_{\xi^N}\right\}. \quad (5.3)$$

The main observation (also used in [PST12]) is that $Q^N(x, \xi^N)$ can be approximated by a Gaussian random variable

$$Q^N(x, \xi^N) \approx Z_{\ell} \quad (5.4)$$

where $Z_{\ell} \overset{\mathcal{D}}{\sim} \mathcal{N}(-\frac{\ell^3}{3}, \frac{\ell^3}{2})$. These approximations are made rigorous in Lemma 7.5 and Lemma 7.6. Therefore, the Bernoulli random variable $\gamma^N(x, \xi^N)$ with success probability $1 \wedge e^{Q^N(x,\xi^N)}$ can be approximated by a Bernoulli random variable, independent of $x$, with success probability equal to

$$\alpha(\ell) = \mathbb{E}[1 \wedge e^{Z_{\ell}}]. \quad (5.5)$$

Thus, the limiting acceptance probability of the MALA algorithm is as given in Equation (5.5). Recall that $\Delta t = N^{-\frac{1}{3}}$. With this notation we introduce the drift function $d^N : \mathcal{H}^s \rightarrow \mathcal{H}^s$ given by

$$d^N(x) = \left(h(\ell)\Delta t\right)^{-1}\mathbb{E}\left[x^{1,N} - x^{0,N} | x^{0,N} = x\right] \quad (5.6)$$

and the martingale difference array $\{\Gamma^{k,N} : k \geq 0\}$ defined by $\Gamma^{k,N} = \Gamma^{N}(x^{k,N}, \xi^{k,N})$ with

$$\Gamma^{k,N} = \left(2h(\ell)\Delta t\right)^{-\frac{1}{2}}\left(x^{k+1,N} - x^{k,N} - h(\ell)\Delta t \ d^N(x^{k,N})\right). \quad (5.7)$$
The normalization constant \( h(\ell) \) defined in Equation (5.1) ensures that the drift function \( d^N \) and the martingale difference array \( \{\Gamma^{k,N}\} \) are asymptotically independent from the parameter \( \ell \). The drift-martingale decomposition of the Markov chain \( \{x^{k,N}\}_k \) then reads

\[
x^{k+1,N} - x^{k,N} = h(\ell)\Delta t d^N(x^{k,N}) + \sqrt{2h(\ell)\Delta t} \Gamma^{k,N}.
\] (5.8)

Lemma 7.8 and Lemma 7.9 exploit the Gaussian behaviour of \( Q^N(x, \xi^N) \) described in Equation (5.4) in order to give quantitative versions of the following approximations,

\[
d^N(x) \approx \mu(x) \quad \text{and} \quad \Gamma^{k,N} \approx N(0,C)
\] (5.9)

where \( \mu(x) = -\left(x + C\nabla\Psi(x)\right) \). From Equation (5.8) it follows that for large \( N \) the evolution of the Markov chain resembles the Euler discretization of the limiting diffusion (1.16). The next step consists of proving an invariance principle for a rescaled version of the martingale difference array \( \{\Gamma^{k,N}\} \). The continuous process \( W^N \in \mathcal{C}([0;T], \mathcal{H}^s) \) is defined as

\[
W^N(t) = \sqrt{\Delta t} \sum_{j=0}^{k} \Gamma^{j,N} + \frac{t - k\Delta t}{\sqrt{\Delta t}} \Gamma^{k+1,N} \quad \text{for} \quad k\Delta t \leq t < (k+1)\Delta t.
\] (5.10)

The sequence of processes \( \{W^N\} \) converges weakly in \( \mathcal{C}([0;T], \mathcal{H}^s) \) to a Brownian motion \( W \) in \( \mathcal{H}^s \) with covariance operator equal to \( C_s \). Indeed, Proposition 7.10 proves the stronger result

\[
(x^{0,N}, W^N) \Rightarrow (z^0, W)
\]

where \( \Rightarrow \) denotes weak convergence in \( \mathcal{H}^s \times \mathcal{C}([0;T], \mathcal{H}^s) \) and \( z^0 \sim \pi \) is independent of the limiting Brownian motion \( W \). Once we have the invariance principle and the converge of the drift and diffusion terms, the “Master Theorem” in [PST12] (see Proposition 3.1 of [PST12]) gives the required diffusion limit.

6. Proof of the Main Result

In this section, we give the proof of the Theorem 5.1. To this end, we use Proposition 3.1 of [PST12]. According to Proposition 3.1 of [PST12], to show the diffusion limit, we must show the following three conditions.

1. **Convergence of initial conditions**: \( \pi^N \) converges in distribution to the probability measure \( \pi \) where \( \pi \) has a finite first moment, that is \( \mathbb{E}^\pi[\|x\|] < \infty \).
2. **Invariance principle**: the sequence \( (x^{0,N}, W^N) \) defined by Equation (5.10) converges weakly in \( \mathcal{H}^s \times \mathcal{C}([0;T], \mathcal{H}^s) \) to \( (z^0, W) \) where \( z^0 \sim \pi \) and \( W \) is a Brownian motion in \( \mathcal{H}^s \), independent from \( z^0 \), with covariance operator \( C_s \).
3. **Convergence of the drift**: there exists a globaly Lipschitz function \( \mu : \mathcal{H}^s \rightarrow \mathcal{H}^s \) that satisfies

\[
\lim_{N \rightarrow \infty} \mathbb{E}^{\pi^N}[\|d^N(x) - \mu(x)\|] = 0.
\]

Item (1.) above follows from Lemma 4.3 of [PST12]); also see Section 7.1 below. Item (2.) is proved in Proposition 7.10. Item (3.) is proved in Lemma 7.8. Thus we have established all three conditions required by Proposition 3.1 of [PST12] and thus the proof of our main result is finished.

7. Key Estimates

In this section, we prove some key estimates for the proximal operator, and and also collect some key approximation properties of \( \mu^N \) and \( \pi^N \) from [PST12]. These properties will be repeatedly used throughout.
7.1. Approximation properties of $\mu^N$ and $\pi^N$

- For $\pi_0$-almost every function $x \in \mathcal{H}^s$, the approximation $\mu^N(x) \approx \mu(x)$ holds as $N$ goes to infinity. Indeed, under Assumption 3.1, the sequences of functions $\mu^N : \mathcal{H}^s \to \mathcal{H}^s$ satisfies (see Lemma 4.1 of [PST12]),

$$\pi_0\left( \left\{ x \in \mathcal{H}^s : \lim_{N \to \infty} \|\mu^N(x) - \mu(x)\|_s = 0 \right\} \right) = 1. \quad (7.1)$$

- Under the Assumptions 3.1 the normalization constants $M_{q,N}$ are uniformly bounded so that for any measurable functional $f : \mathcal{H} \mapsto \mathbb{R}$, we have from Lemma 4.3 of [PST12] that

$$\mathbb{E}^{\pi^N}[|f(x)|] \lesssim \mathbb{E}^{\pi_0}[|f(x)|].$$

Moreover, the sequence of probability measure $\pi^N$ satisfies $\pi^N \Rightarrow \pi$ where $\Rightarrow$ denotes weak convergence in $\mathcal{H}^s$.

- Fernique’s theorem [DPZ92] states that for any exponent $p \geq 0$ we have $\mathbb{E}^{\pi^N}[\|x\|_p^p] < \infty$. We also have that for any $p \geq 0$

$$\sup_{N \in \mathbb{N}} \mathbb{E}^{\pi^N}[\|x\|_p^p] < \infty.$$

7.2. Estimates involving proximal functions and the remainder term

Recall the constant $M_6$ from (3.14).

Lemma 7.1. For any $x \in \mathcal{H}^s$ and $N \in \mathbb{N}$ and for all $\delta < \frac{1}{2M_6}$,

$$\|\text{Prox}_{\delta,N}(x) - x\|_s \lesssim \delta (1 + \|x\|_s).$$

Proof. Set $x^* = \text{Prox}_{\delta,N}(x)$. Since $x^*$ minimizes the map:

$$y \mapsto (\Psi^N(y) + \frac{1}{2\delta} \|y - x\|_s^2),$$

from our assumptions in (3.14) and (3.7), it follows that

$$\frac{1}{2\delta} \|x^* - x\|_s^2 \leq \Psi^N(x) - \Psi^N(x^*) = |\Psi^N(x^*) - \Psi^N(x)|$$

$$\leq |\langle \nabla \Psi^N(x), x^* - x \rangle| + M_6 \|x^* - x\|_s^2$$

$$\leq M_3 (1 + \|x\|_s) \|x^* - x\|_s + M_6 \|x^* - x\|_s^2.$$

Dividing by the term $\|x^* - x\|_s$ throughout and simplifying yields

$$\|x^* - x\|_s \leq \delta \frac{M_3}{1 - 2\delta M_6} (1 + \|x\|_s) \lesssim \delta (1 + \|x\|_s)$$

and the proof is done. \qed

Lemma 7.2. Recall the remainder term $R^N(x, \delta)$ from (4.8). For any $x \in \mathcal{H}^s$, $N \in \mathbb{N}$ and for all $\delta < \frac{1}{2M_6}$,

$$\|R^N(x, \delta)\|_{\mathcal{C}^N} \lesssim \delta^2 (1 + \|x\|_s), \quad \|R^N(x, \delta)\|_s \lesssim \delta^2 (1 + \|x\|_s).$$

Proof. Set $x^* = \text{Prox}_{\delta,N}(x)$. Then $R^N(x, \delta) = \delta \mathcal{C}^N(x^* - x)$. Thus

$$\|R^N(x, \delta)\|_{\mathcal{C}^N}^2 = \langle R^N(x, \delta), (\mathcal{C}^N)^{-1} R^N(x, \delta) \rangle$$

$$= \delta^2 \langle \mathcal{C}^N(x^* - x), (x^* - x) \rangle$$
where the last inequality follows from Lemma 7.1 showing the first inequality. The second inequality follows similarly:
\[ \| R^N(x, \delta) \|^2 = \delta^2 \| C^N(x^* - x) \|^2 \lesssim \delta^2 \| x^* - x \|^2 \lesssim \delta^4 (1 + \| x \|^2) \]
and the proof is done. \( \square \)

Next lemma shows that the size of the jump \( y - x \) is of order \( \sqrt{\Delta t} \).

**Lemma 7.3.** Consider \( y \) given by (4.9). Under Assumptions 3.1, for any \( p \geq 1 \) we have
\[ \mathbb{E}^{x, N}_x [ \| y - x \|^p ] \lesssim (\Delta t)^{\frac{p}{2}} \cdot (1 + \| x \|^p) \].

**Proof.** Under Assumption 3.1 the function \( \mu^N \) is globally Lipschitz on \( \mathcal{H}^s \), with Lipschitz constant that can be chosen independent from \( N \). Thus using Lemma 7.2 we obtain that
\[ \| y - x \|_s \lesssim \Delta t (1 + \| x \|_s) + \| R^N(x, \delta) \|_s + \sqrt{\Delta t} \| C^N \xi \|_s \]
\[ \lesssim \Delta t (1 + \| x \|_s) + (\Delta t)^2 (1 + \| x \|_s) + \sqrt{\Delta t} \| C^N \xi \|_s \]
\[ \lesssim \Delta t (1 + \| x \|_s) + \sqrt{\Delta t} \| C^N \xi \|_s \].
We have \( \mathbb{E}^{x, 0}[\| C^N \xi \|_s] \leq \mathbb{E}^{x, 0}[\| \xi \|_s] < \infty \), where \( \xi \overset{D}{\sim} N(0, \mathcal{C}) \). Consequently, \( \mathbb{E}^{x, 0}[\| C^N \xi \|_s] \) is uniformly bounded as a function of \( N \), proving the lemma. \( \square \)

Consider \( y \) given by (4.9) and recall from (4.7) that
\[ y = x_{\text{MALA}} + R^N(x, \delta) \].

**Lemma 7.4.** We have
\[ a^N(x, \delta) \equiv \| y \|_{C^N}^2 - \| x_{\text{MALA}} \|_{C^N}^2 \]
\[ \mathbb{E}^{x, N}_x a^N(x, \delta) \lesssim \delta^2 \]

**Proof.** From (4.7) we have
\[ \| y \|_{C^N}^2 - \| x_{\text{MALA}} \|_{C^N}^2 = a^N(x, \delta) \]
\[ a^N(x, \delta) \equiv 2(x_{\text{MALA}}, R^N(x))_{C^N} + \| R^N(x, \delta) \|_{C^N}^2. \] (7.2)

From (4.8), we obtain
\[ |\langle x_{\text{MALA}}, R^N(x, \delta) \rangle_{C^N}| = |\langle x_{\text{MALA}}, (C^N)^{-1} R^N(x, \delta) \rangle| \]
\[ \leq \| x_{\text{MALA}} \|_s \| R^N(x, \delta) \|_{C^N} \].
From Lemma 7.3 we deduce that
\[ \| x_{\text{MALA}} \|_s \lesssim (1 + \| x \|_s) + \sqrt{\Delta t} \| C^N \xi \|_s \].
Combining this with Lemma 7.2 yields that
\[ |\langle x_{\text{MALA}}, R^N(x, \delta) \rangle_{C^N}| \lesssim \delta^2 (1 + \| x \|_s^2) (1 + \sqrt{\Delta t} \| C^N \xi \|_s) \].

Thus
\[ \mathbb{E}^{x, N}_x (|\langle x_{\text{MALA}}, R^N(x, \delta) \rangle_{C^N}|) \lesssim \delta^2 \mathbb{E}^{x, N}_x (1 + \| x \|_s^2) (1 + \sqrt{\Delta t} \| C^N \xi \|_s) \lesssim \delta^2. \] (7.3)
Thus from (7.2), (7.3) and Lemma 7.2 we deduce that
\[ \mathbb{E}^{x, N}_x (a^N(x, \delta) \lesssim \delta^2 \]
and the proof is finished. \( \square \)
7.3. Gaussian approximation of $Q^N$

Recall the quantity $Q^N$ defined in Equation (5.3). This section proves that $Q^N$ has a Gaussian behavior in the sense that

$$Q^N(x, \xi^N) = Z^N(x, \xi^N) + i^N(x, \xi^N) + e^N(x, \xi^N)$$

(7.4)

where the quantities $Z^N$ and $i^N$ are equal to

$$Z^N(x, \xi^N) = -\frac{\ell^3}{4} + \frac{\ell^2}{\sqrt{2}} N^{-\frac{1}{2}} \sum_{j=1}^{N} \lambda_j^{-1} \xi_j$$

(7.5)

$$i^N(x, \xi^N) = \frac{1}{2} (\ell \Delta t)^2 \left( \| x \|_{2N}^2 - \| (C^N)^N \xi^N \|_{C^N}^2 \right)$$

(7.6)

with $i^N$ and $e^N$ small. Thus the principal contributions to $Q^N$ comes from the random variable $Z^N(x, \xi^N)$. Notice that, for each fixed $x \in \mathcal{H}^s$, the random variable $Z^N(x, \xi^N)$ is Gaussian. Furthermore, the Karhunen-Loève expansion of $\pi_0$ shows that for $\pi_0$-almost every choice of function $x \in \mathcal{H}$ the sequence $\{Z^N(x, \xi^N)\}_{N \geq 1}$ converges in law to the distribution of $Z_\ell \mathcal{N}(0, \frac{\ell^2}{\ell^2 + N})$. The next lemma rigorously bounds the error terms $e^N(x, \xi^N)$ and $i^N(x, \xi^N)$: we show that $i^N$ is an error term of order $O(N^{-\frac{1}{2}})$ and $e^N(x, \xi)$ is an error term of order $O(\sqrt{N})$. In Lemma 7.6 we then quantify the convergence of $Z^N(x, \xi^N)$ to $Z_\ell$.

**Lemma 7.5. (Gaussian Approximation)** Let $p \geq 1$ be an integer. Under Assumptions 3.1, $Q^N(x, \xi^N)$ has the expansion given in (7.4) and the error terms $i^N$ and $e^N$ in the Gaussian approximation (7.4) satisfy

$$\left( \mathbb{E}^N [i^N(x, \xi^N)|p] \right)^{\frac{1}{2}} = O(N^{-\frac{1}{2}}) \quad \text{and} \quad \left( \mathbb{E}^N [e^N(x, \xi^N)|p] \right)^{\frac{1}{2}} = O(N^{-\frac{1}{2}}).$$

(7.7)

**Proof.** As in Lemma 4.4 of [PST12], without loss of generality, we suppose $p = 2q$. The quantity $Q^N$ is defined in Equation (5.3) and expanding terms leads to

$$Q^N(x, \xi^N) = I_1 + I_2 + I_3 + I_4$$

where the quantities $I_1$, $I_2$, $I_3$ and $I_4$ are given by

$$I_1 = -\frac{1}{2} \left( \| y \|_{2N}^2 - \| x \|_{2N}^2 \right) - \frac{1}{4 \ell \Delta t} \left( \| x - y(1 - \ell \Delta t) \|_{2N}^2 - \| y - x(1 - \ell \Delta t) \|_{2N}^2 \right)$$

$$I_2 = -\left( \Psi^N(y) - \Psi^N(x) \right) - \frac{1}{2} \left( (x - y(1 - \ell \Delta t), C^N \nabla \Psi^N(y))_{CN} - (y - x(1 - \ell \Delta t), C^N \nabla \Psi^N(x))_{CN} \right)$$

$$I_3 = -\frac{1}{4 \ell \Delta t} \left\{ \| \ell \Delta t C^N \nabla \Psi^N(y) + R^N(y, \delta) \|_{2N}^2 - \| \ell \Delta t C^N \nabla \Psi^N(x) + R^N(x, \delta) \|_{2N}^2 \right\}$$

$$I_4 = -\frac{1}{2 \ell \Delta t} \left\{ (x - y(1 - \ell \Delta t), R^N(y, \delta))_{CN} - (y - x(1 - \ell \Delta t), R^N(x, \delta))_{CN} \right\}.$$

The term $I_1$ arises purely from the Gaussian part of the target measure $\pi^N$ and from the Gaussian part of the proposal. The other terms come from the change of probability involving the functional $\Psi^N$. By the calculation identical to page 2343 of [PST12], we can simplify the the term $I_1$ to be:

$$I_1 = -\frac{\ell \Delta t}{4} \left( \| y \|_{2N}^2 - \| x \|_{2N}^2 \right).$$

(7.8)

The term $I_1$ is shown to be $O(1)$ and constitutes the main contribution to $Q^N$. Before analyzing $I_1$ in more detail, we show that $I_2$, $I_3$ and $I_4$ are $O(N^{-\frac{1}{2}})$:

$$\left( \mathbb{E}^N [I_2^2] \right)^{\frac{1}{2}} + \left( \mathbb{E}^N [I_3^2] \right)^{\frac{1}{2}} + \left( \mathbb{E}^N [I_4^2] \right)^{\frac{1}{2}} = O(N^{-\frac{1}{2}}).$$

(7.9)
• By a calculation nearly identical to the one in Lemma 4.4 of [PST12] (the only change being the use of our Lemma 7.3 instead of their Lemma 4.2) we obtain that
\[
\left( E^{π^N} [I^2_q] \right)^{1/2} = O(N^{-\frac{1}{2}}). \tag{7.10}
\]

• Using the definition of \( R^N(x, \delta) \) from (4.8), we obtain that
\[
\begin{align*}
E^{π^N} [I^2_q] &\lesssim Δt^{2q} E^{π^N} \left[ |⟨\nabla Ψ^N(x), C^N \nabla Ψ^N(x)⟩|^q + |⟨\nabla Ψ^N(y), C^N \nabla Ψ^N(y)⟩|^q \right] \\
&\quad + Δt^{-2q} E^{π^N} \left[ ∥R^N(x, \delta)∥_{C^N}^2 + ∥R^N(y, \delta)∥_{C^N}^2 \right].
\end{align*}
\]

Lemma 3.5 states \( C^N \nabla Ψ^N : \mathcal{H}^s \to \mathcal{H}^s \) is globally Lipschitz, with a Lipschitz constant that can be chosen uniformly in \( N \). Therefore,
\[
∥C^N \nabla Ψ^N(z)∥_s \lesssim 1 + ∥z∥_s. \tag{7.11}
\]

Since \( ∥C^N \nabla Ψ^N(z)∥_{C^N}^2 = ⟨\nabla Ψ^N(z), C^N \nabla Ψ^N(z)⟩ \), the bound (3.7) gives
\[
E^{π^N} [I^2_q] \lesssim Δt^{2q} E^{π^N} \left[ ⟨\nabla Ψ^N(x), C^N \nabla Ψ^N(x)⟩^q + ⟨\nabla Ψ^N(y), C^N \nabla Ψ^N(y)⟩^q \right]
\lesssim Δt^{2q} E^{π^N} \left[ (1 + ∥x∥_s)^{2q} + (1 + ∥y∥_s)^{2q} \right]
\lesssim Δt^{2q} E^{π^N} \left[ 1 + ∥x∥_s^{2q} + ∥y∥_s^{2q} \right] \lesssim Δt^{2q} = \left( N^{-\frac{1}{2}} \right)^{2q}. \tag{7.12}
\]

Similarly, from Lemma 7.2 and 7.3,
\[
Δt^{-2q} E^{π^N} \left[ ∥R^N(x, \delta)∥_{C^N}^2 + ∥R^N(y, \delta)∥_{C^N}^2 \right] \lesssim Δt^{6q} E^{π^N} \left[ 1 + ∥x∥_s^{2q} + ∥y∥_s^{2q} \right] \lesssim Δt^{6q}
\lesssim \left( N^{-\frac{1}{2}} \right)^{6q} \lesssim \left( N^{-\frac{1}{2}} \right)^{2q}. \tag{7.13}
\]

Thus from (7.12) and (7.13), we conclude that
\[
\left( E^{π^N} [I^2_q] \right)^{1/2} \lesssim \left( N^{-\frac{1}{2}} \right)^{2q}. \tag{7.14}
\]

• Finally, we tackle the term \( I_1 \):
\[
E^{π^N} [I^2_q] \lesssim Δt^{-2q} E^{π^N} \left[ ∥x - y(1 - ℓΔt)∥_{s}^{2q} ∥(C^N)^{-1} R^N(y, \delta)∥_{s}^{2q}
\quad + ∥y - x(1 - ℓΔt)∥_{s}^{2q} ∥(C^N)^{-1} R^N(x, \delta)∥_{s}^{2q} \right].
\]

From Lemma 7.3, we obtain that \( E^{π^N} (∥y - x(1 - ℓΔt)∥_{s}^{4q}) \lesssim (Δt)^{2q} · (1 + ∥x∥_{s}^{4q}) \) and \( E^{π^N} (∥x - y(1 - ℓΔt)∥_{s}^{4q}) \lesssim (Δt)^{2q} · (1 + ∥x∥_{s}^{4q}). \) Similarly, from Lemma 7.2 we gather that \( E^{π^N} ∥R^N(x, \delta)∥_{C^N}^{4q} \lesssim δ^{8q} (1 + ∥x∥_{s}^{4q}). \) Putting these two together and using the Cauchy-Schwartz inequality gives,
\[
E^{π^N} [I^2_q] \lesssim Δt^{3q} E^{π^N} \left[ 1 + ∥x∥_s^{2q} + ∥y∥_s^{2q} \right] \lesssim \left( N^{-\frac{1}{2}} \right)^{2q}. \tag{7.15}
\]

Equations (7.10), (7.14) and (7.15) imply the requisite estimate in (7.9).

Next, we tackle the term \( I_1 \). Recall from from (7.8) that
\[
I_1 = -\frac{ℓΔt}{4} ∥y∥_{C^N}^2 - ∥x∥_{C^N}^2.
\]

From Lemma 7.4 we obtain that
\[
∥y∥_{C^N}^2 = ∥x_{MALA}∥_{C^N}^2 + a^N(x, ℓΔt), \quad E^{π^N} a^N(x, ℓΔt) \lesssim (Δt)^2. \tag{7.16}
\]
Consequently,

\[ I_3 = -\frac{\Delta t}{4} \left( \|x_{\text{MALA}}\|^2_{CN} - \|x\|^2_{CN} \right) - \frac{\Delta t}{4} \alpha(x, \ell \Delta t). \]

From Lemma 4.4 of [PST12], we deduce that

\[ -\frac{\Delta t}{4} \left( \|x_{\text{MALA}}\|^2_{CN} - \|x\|^2_{CN} \right) = Z_N(x, \xi_N) + i_N(x, \xi_N) + b_N(x, \xi_N) \]

with \( Z_N(x, \xi_N) \) and \( i_N(x, \xi_N) \) given by Equation (7.5) and (7.6) and

\[ \left( \mathbb{E}^\pi \left[ b_N(x, \xi_N)^2 \right] \right)^{\frac{1}{2}} = O(N^{-\frac{1}{2}}). \]

Lemma 4.4 of [PST12] also shows that

\[ \left( \mathbb{E}^\pi \left[ i_N(x, \xi_N)^2 \right] \right)^{\frac{1}{2}} = O(N^{-\frac{1}{2}}). \]

The proof of the lemma now follows from (7.9), (7.16) and (7.17). \[ \square \]

We recall Lemma 4.5 of [PST12]:

**Lemma 7.6.** [PST12, Lemma 4.5] (Asymptotic independence) Let \( p \geq 1 \) be a positive integer and \( f : \mathbb{R} \to \mathbb{R} \) be a 1-Lipschitz function. Consider error terms \( e_N(x, \xi) \) satisfying

\[ \lim_{N \to \infty} \mathbb{E}^\pi[e_N^p(x, \xi_N)] = 0. \]

Define the functions \( f^N : \mathbb{R} \to \mathbb{R} \) and the constant \( f \in \mathbb{R} \) by

\[ f^N(x) = \mathbb{E}_x [f(Z(x, \xi_N) + e^*_N(x, \xi_N))] \quad \text{and} \quad f = \mathbb{E}[f(Z)]. \]

Then the function \( f^N \) is highly concentrated around its mean in the sense that

\[ \lim_{N \to \infty} \mathbb{E}^\pi \left[ |f^N - f|^p \right] = 0. \]

**Corollary 7.7.** Let \( p \geq 1 \) be a positive. The local mean acceptance probability \( \alpha^N(x) \) defined in Equation (4.11) satisfies

\[ \lim_{N \to \infty} \mathbb{E}^\pi \left[ |\alpha^N(x) - \alpha(\ell)|^p \right] = 0. \]

**Proof.** The function \( f(z) = 1 \wedge e^z \) is 1-Lipschitz and \( \alpha(\ell) = \mathbb{E}[f(Z)]. \) Also,

\[ \alpha^N(x) = \mathbb{E}_x [f(Q_N(x, \xi_N))] = \mathbb{E}_x [f(Z_N(x, \xi_N) + e^*_N(x, \xi_N))] \]

with \( e^*_N(x, \xi_N) = i_N(x, \xi_N) + e_N(x, \xi_N). \) Lemma 7.5 shows that \( \lim_{N \to \infty} \mathbb{E}^\pi \left[ e^*_N(x, \xi)^p \right] = 0 \) and therefore Lemma 7.6 gives the conclusion. \[ \square \]

### 7.4. Drift approximation

This section proves that the approximate drift function \( d^N : \mathcal{H}^s \to \mathcal{H}^s \) defined in Equation (5.6) converges to the drift function \( \mu : \mathcal{H}^s \to \mathcal{H}^s \) of the limiting diffusion (5.2).

**Lemma 7.8.** (Drift Approximation) Let Assumptions 3.1 hold. The drift function \( d^N : \mathcal{H}^s \to \mathcal{H}^s \) converges to \( \mu \) in the sense that

\[ \lim_{N \to \infty} \mathbb{E}^\pi \left[ \|d^N(x) - \mu(x)\|^2 \right] = 0. \]
Proof. Now that we have established the relevant estimates, the proof of this lemma is nearly identical to that of Lemma 4.7 of [PST12], but also needs to account for the extra error term induced by the proximal operator. The approximate drift \( d^N \) is given by Equation (5.6). The definition of the local mean acceptance probability \( \alpha^N(x) \) given by Equation (4.11) shows that \( d^N \) can also be expressed as

\[
d^N(x) = \left( \alpha^N(x) \alpha(\ell)^{-1} \right) \mu^N(x) + \mathcal{R}^N_{\text{Prox}}(x, \Delta t) + \sqrt{2e}h(\ell)^{-1}(\Delta t)^{-\frac{4}{3}}\varepsilon^N(x) \tag{7.19}
\]

where \( \mu^N(x) = -(P^N x + C^N \nabla \Psi^N(x)) \); the term \( \varepsilon^N(x) \) is defined by

\[
\varepsilon^N(x) = E_x [\gamma^N(x, \xi^N)C^\dagger \xi^N] = E_x \left[ (1 \wedge e^{Q^N(x, \xi^N)}) C^\dagger \xi^N \right]
\]

and the term \( \mathcal{R}^N_{\text{Prox}}(x, \Delta t) \) is the error term induced by the proximal approxiation:

\[
\mathcal{R}^N_{\text{Prox}}(x, \Delta t) = \frac{\alpha^N(x) R^N(x, \ell \Delta t)}{h(\ell)} \Delta t.
\]

To prove Lemma 7.8 it suffices to verify that

\[
\lim_{N \to \infty} \mathbb{E}^{\pi^N} \left[ \left| \left( \alpha^N(x) \alpha(\ell)^{-1} \right) \mu^N(x) - \mu(x) \right|^2 \right] = 0 \tag{7.20}
\]

\[
\lim_{N \to \infty} \mathbb{E}^{\pi^N} \left[ \left| \mathcal{R}^N_{\text{Prox}}(x, \Delta t) \right|^2 \right] = 0 \tag{7.21}
\]

\[
\lim_{N \to \infty} (\Delta t)^{-1} \mathbb{E}^{\pi^N} \left[ \left| \varepsilon^N(x) \right|^2 \right] = 0 \tag{7.22}
\]

- Equation (7.20) follows directly from Lemma 4.7 of [PST12].
- Next, using the fact that \(|\alpha^N(x)| \leq 1 \) and Lemma 7.2,

\[
\left| \mathcal{R}^N_{\text{Prox}}(x, \Delta t) \right|^2 \lesssim \left( \alpha^N(x) \right)^2 \left| \frac{R^N(x, \ell \Delta t)}{\Delta t} \right|^2 \lesssim (\Delta t)^2 (1 + \|x\|^2_2)
\]

and thus we have

\[
\lim_{N \to \infty} \mathbb{E}^{\pi^N} \left| \mathcal{R}^N_{\text{Prox}}(x, \Delta t) \right|^2 = \lim_{N \to \infty} N^{-2/3} \mathbb{E}^{\pi^N} (1 + \|x\|^2_2) = 0
\]

establishing (7.21).
- Let us prove Equation (7.22). If the Bernoulli random variable \( \gamma^N(x, \xi^N) \) were independent from the noise term \((C^N)^\dagger \xi^N \), it would follow that \( \varepsilon^N(x) = 0 \). In general \( \gamma^N(x, \xi^N) \) is not independent from \((C^N)^\dagger \xi^N \) so that \( \varepsilon^N(x) \) is not equal to zero. Nevertheless, as quantified by Lemma 7.6, the Bernoulli random variable \( \gamma^N(x, \xi^N) \) is asymptotically independent from the current position \( x \) and from the noise term \((C^N)^\dagger \xi^N \). Consequently, we can prove in Equation (7.24) that the quantity \( \varepsilon^N(x) \) is small. To this end, we establish that each component \( \langle \varepsilon(x), \hat{\varphi}_j \rangle_i^2 \) satisfies

\[
\mathbb{E}^{\pi^N} \left[ \langle \varepsilon^N(x), \hat{\varphi}_j \rangle_i^2 \right] \lesssim N^{-1} \mathbb{E}^{\pi^N} \left[ \langle x, \hat{\varphi}_j \rangle_i^2 \right] + N^{-\frac{1}{2}} (\lambda_j^2) N^{-2} \mathbb{E}^{\pi^N} \left[ \langle x, \hat{\varphi}_j \rangle_i^2 \right]. \tag{7.23}
\]

Summation of Equation (7.23) over \( j = 1, \ldots, N \) leads to

\[
\mathbb{E}^{\pi^N} \left[ \|\varepsilon^N(x)\|^2 \right] \lesssim N^{-1} \mathbb{E}^{\pi^N} \left[ \|x\|^2 \right] + N^{-\frac{1}{2}} \text{Tr}_{\mathcal{H}^N}(\mathcal{C}_x) \lesssim N^{-\frac{1}{4}}, \tag{7.24}
\]

which gives the proof of Equation (7.22). To prove Equation (7.23) for a fixed index \( j \in \mathbb{N} \), the quantity \( Q^N(x, \xi) \) is decomposed as a sum of a term independent from \( \xi_j \) and another remaining term of small magnitude. To this end we introduce

\[
\begin{align*}
Q_j^N(x, \xi^N) &= Q_j^N(x, \xi^N) + Q_j^{N, \perp}(x, \xi^N) \\
Q_j^N(x, \xi^N) &= -\frac{1}{\sqrt{2}} \ell^2 \lambda_j^2 \xi_j - \frac{1}{2} \ell^2 N^{-\frac{1}{2}} \lambda_j^2 \xi_j^2 + e^N(x, \xi^N). \tag{7.25}
\end{align*}
\]
The definitions of $Z^N(x, \xi^N)$ and $i^N(x, \xi^N)$ in Equation (7.5) and (7.6) readily show that $Q^N_{j,\pm}(x, \xi^N)$ is independent from $\xi_j$. The noise term satisfies $C^* \xi^N = \sum_{j=1}^{N} (j^* \lambda_j) \xi_j \hat{\varphi}_j$. Since $Q^N_{j,\pm}(x, \xi^N)$ and $\xi_j$ are independent and $z \mapsto 1 \wedge e^z$ is 1-Lipschitz, it follows that

$$
\langle \varepsilon^N(x), \hat{\varphi}_j \rangle^2 = (j^* \lambda_j)^2 \left( \mathbb{E}_x \left[ (1 \wedge e^{Q^N(x, \xi^N)}) \xi_j \right] \right)^2 \\
= (j^* \lambda_j)^2 \left( \mathbb{E}_x \left[ (1 \wedge e^{Q^N(x, \xi^N)}) - (1 \wedge e^{Q^N_{j,\pm}(x, \xi^N)}) \right] \xi_j \right)^2 \\
\lesssim (j^* \lambda_j)^2 \mathbb{E}_x \left[ (Q^N(x, \xi^N)) - Q^N_{j,\pm}(x, \xi^N) \right] \\
= (j^* \lambda_j)^2 \mathbb{E}_x \left[ Q^N_{j,\pm}(x, \xi^N) \right]^2.
$$

By Lemma 7.5 $\mathbb{E}^N [e^N(x, \xi^N)^2] \lesssim N^{-\frac{2}{3}}$. Therefore,

$$
(j^* \lambda_j)^2 \mathbb{E}^N [Q^N_{j,\pm}(x, \xi^N)^2] \lesssim (j^* \lambda_j)^2 \left\{ N^{-1} \lambda_j^{-2} \mathbb{E}^N [x_j^2 \xi_j^2] + N^{-\frac{4}{3}} \mathbb{E}^N [\lambda_j^2 \xi_j^4] + \mathbb{E}^N [e^N(x, \xi)^2] \right\} \\
\lesssim N^{-1} \mathbb{E}^N [(j^* x_j)^2 \xi_j^2] + (j^* \lambda_j)^2 (N^{-\frac{4}{3}} + N^{-\frac{2}{3}}) \\
\lesssim N^{-1} \mathbb{E}^N [(x, \hat{\varphi}_j)^2] + (j^* \lambda_j)^2 N^{-\frac{2}{3}} \\
\lesssim N^{-1} \mathbb{E}^N [(x, \hat{\varphi}_j)^2] + (j^* \lambda_j)^2 N^{-\frac{2}{3}},
$$

which finishes the proof of Equation (7.23).

Thus we have established (7.20), (7.21) and (7.22) and the proof is finished. \hfill \square

### 7.5. Noise approximation

Recall the definition (5.7) of the martingale difference $\Gamma^{k,N}$. In this section we estimate the error in the approximation $\Gamma^{k,N} \approx N(0, \mathcal{C}_s)$. To this end we introduce the covariance operator

$$
D^N(x) = \mathbb{E}_x \left[ \Gamma^{k,N} \otimes \mathcal{H}, \Gamma^{k,N} | x^{k,N} = x \right].
$$

For any $x, u, v \in \mathcal{H}^s$ the operator $D^N(x)$ satisfies

$$
\mathbb{E} \left[ (\Gamma^{k,N}, u)_s (\Gamma^{k,N}, v)_s | x^{k,N} = x \right] = \langle u, D^N(x) v \rangle_s.
$$

The next lemma gives a quantitative version of the approximation $D^N(x) \approx \mathcal{C}_s$.

**Lemma 7.9.** Let Assumptions 3.1 hold. For any pair of indices $i, j \geq 0$ the operator $D^N(x) : \mathcal{H}^s \to \mathcal{H}^s$ satisfies

$$
\lim_{N \to \infty} \mathbb{E}^N \left| \langle \hat{\varphi}_i, D^N(x) \hat{\varphi}_j \rangle_s - \langle \hat{\varphi}_i, C_s \hat{\varphi}_j \rangle_s \right| = 0 \tag{7.26}
$$

and, furthermore,

$$
\lim_{N \to \infty} \mathbb{E}^N \left| \text{Tr}_{\mathcal{H}^s} (D^N(x)) - \text{Tr}_{\mathcal{H}^s} (C_s) \right| = 0. \tag{7.27}
$$

**Proof.** This lemma follows directly from Lemma 4.8 of [PST12], since the only estimate needed for the proof of Lemma 4.8 of [PST12] is the Gaussian approximation and the estimate for $e^N(x, \xi^N)$ established in Lemma 7.5. Thus the proof is finished. \hfill \square

### 7.6. Martingale Invariance Principle

This section proves that the process $W^N$ defined in Equation (5.10) converges to a Brownian motion.
Proposition 7.10. Let Assumptions 3.1 hold. Let $z^0 \sim \pi$ and $W^N(t)$ the process defined in equation (5.10) and $x^{0,N} \overset{D}{=} \pi^N$ the starting position of the Markov chain $x^N$. Then

$$(x^{0,N}, W^N) \implies (z^0, W), \quad (7.28)$$

where $\implies$ denotes weak convergence in $\mathcal{H}^s \times C([0, T]; \mathcal{H}^s)$, and $W$ is a $\mathcal{H}^s$-valued Brownian motion with covariance operator $C_s$. Furthermore the limiting Brownian motion $W$ is independent of the initial condition $z^0$.

Proof. This proof involves verifying three conditions of Proposition 5.1 of [Ber86] and is identical to that of Proposition 4.10 of [PST12]. The only change required is to use our Lemma 7.9 instead of their Lemma 4.8. Therefore we omit the details of the rest of the proof.

8. Closing Comments

There are a number of issues that could be followed up in future work that are of great practical interest:

- A straightforward, but tedious extension will be to extend Theorem 2.1 to a general class of product measures.
- As mentioned in the introduction, we only study the case when the log-target is differentiable. Of course, the most interesting case is when the log-target is not differentiable. In this context, it would be very worthwhile to understand how the scaling differs from the usual MALA for non-smooth target distributions.
- We also set $\lambda = \delta$; it is not clear to us if this is necessary. Are there regimes when $\lambda$ and $\delta$ scaled differently as a function of the dimension $N$ that are better than $\lambda = \delta$?
- A similar result should be of interest when proximal functions are used for implementing the Hybrid Monte Carlo algorithm.

Acknowledgements. The author thanks Gareth Roberts and Andrew Stuart for introducing him to optimal scaling of Markov Chain Monte Carlo Methods and several inspiring conversations, Dean Foster and Robert Stine for introducing him to proximal algorithms and related ADMM methods, Alex Belloni, David Dunson, Lucas Janson, Eric Laber, Jonathan Mattingly, Sayan Mukherjee, Debdeep Pati, Christian Robert, Aaron Smith, Pragya Sur and Alex Thiery for very helpful comments and their collaboration over the years. Part of this work was done when the author was an Amazon Scholar. He thanks Robert Stine and Andrew Rund for their encouragement to work on this problem.

References

[BC11] Heinz H Bauschke and Patrick L Combettes. *Convex analysis and monotone operator theory in Hilbert spaces*, volume 408. Springer, 2011.
[Ber86] E. Berger. Asymptotic behaviour of a class of stochastic approximation procedures. *Probab. Theory Relat. Fields*, 71(4):517–552, 1986.
[BRS09] A. Beskos, G.O. Roberts, and A.M. Stuart. Optimal scalings of metropolis-hastings algorithms for non-product targets in high dimensions. *Annals of Applied Probability*, 19:863–898, 2009.
[BRSV08] A. Beskos, G.O. Roberts, A.M. Stuart, and J. Voss. An MCMC method for diffusion bridges. *Stochastics and Dynamics*, 8(3):319–350, 2008.
[BS09] A. Beskos and A.M. Stuart. MCMC methods for sampling function space. In *Invited Lectures, Sixth International Congress on Industrial and Applied Mathematics, ICIAM07*, Editors Rolf Jeltsch and Gerhard Wanner, pages 337–364. European Mathematical Society, 2009.
[DPZ92] G. Da Prato and J. Zabczyk. *Stochastic equations in infinite dimensions*, volume 44 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, Cambridge, 1992.
[HAVW05] M. Hairer, A.M. Stuart, J. Voss, and P. Wiberg. Analysis of SPDEs arising in path sampling. Part I: the gaussian case. *Comm. Math. Sci.*, 3:587–603, 2005.
[HSV07] M. Hairer, A. M. Stuart, and J. Voss. Analysis of SPDEs arising in path sampling. Part II: the nonlinear case. *Ann. Appl. Probab.*, 17(5-6):1657–1706, 2007.

[HSV11] M. Hairer, A. M. Stuart, and J. Voss. Signal processing problems on function space: Bayesian formulation, stochastic PDEs and effective MCMC methods. In D. Crisan and B. Rozovsky, editors, *The Oxford Handbook of Nonlinear Filtering*. Oxford University Press, 2011.

[Lam21] Andrew Lamperski. Projected stochastic gradient langevin algorithms for constrained sampling and non-convex learning. In *Conference on Learning Theory*, pages 2891–2937. PMLR, 2021.

[MPS12] Jonathan C Mattingly, Natesh S Pillai, and Andrew M Stuart. Diffusion limits of the random walk metropolis algorithm in high dimensions. *The Annals of Applied Probability*, 22(3):881–930, 2012.

[Per16] Marcelo Pereyra. Proximal markov chain monte carlo algorithms. *Statistics and Computing*, 26(4):745–760, 2016.

[PST12] Natesh S Pillai, Andrew M Stuart, and Alexandre H Thiéry. Optimal scaling and diffusion limits for the langevin algorithm in high dimensions. *The Annals of Applied Probability*, 22(6):2320–2356, 2012.

[RC04] C. P. Robert and G. Casella. *Monte Carlo statistical methods*. Springer Texts in Statistics. Springer-Verlag, New York, second edition, 2004.

[RGG97] G. O. Roberts, A. Gelman, and W. R. Gilks. Weak convergence and optimal scaling of random walk Metropolis algorithms. *Ann. Appl. Probab.*, 7(1):110–120, 1997.

[RR98] G. O. Roberts and J. S. Rosenthal. Optimal scaling of discrete approximations to Langevin diffusions. *J. R. Stat. Soc. Ser. B Stat. Methodol.*, 60(1):255–268, 1998.

[RR01] G. O. Roberts and J. S. Rosenthal. Optimal scaling for various Metropolis-Hastings algorithms. *Statist. Sci.*, 16(4):351–367, 2001.

[Stu10] A.M. Stuart. Inverse problems: a Bayesian perspective. *Acta Numerica*, 19, 2010.

[WT11] Max Welling and Yee W Teh. Bayesian learning via stochastic gradient langevin dynamics. In *Proceedings of the 28th international conference on machine learning (ICML-11)*, pages 681–688. Citeseer, 2011.