GENERALIZING HAMILTONIAN MONTE CARLO WITH NEURAL NETWORKS

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ABSTRACT

We present a general-purpose method to train Markov Chain Monte Carlo kernels (parameterized by deep neural networks) that converge and mix quickly to their target distribution. Our method generalizes Hamiltonian Monte Carlo and is trained to maximize expected squared jumped distance, a proxy for mixing speed. We demonstrate significant empirical gains (up to $50\times$ greater effective sample size) on a collection of simple but challenging distributions. Finally, we show quantitative and qualitative gains on a real-world task: latent-variable generative modeling. Python source code will be open-sourced with the camera-ready paper.

1 INTRODUCTION

High-dimensional distributions that are only analytically tractable up to a normalizing constant are ubiquitous in many fields. For instance, they arise in protein folding (Schütte et al., 1999), physics simulations (Olsson 1995), and machine learning (Andrieu et al., 2003). Sampling from such distributions is a critical task for learning and inference (MacKay 2003), however it is an extremely hard problem in general.

Markov Chain Monte Carlo (MCMC) methods promise a solution to this problem. They operate by generating a sequence of correlated samples that converge in distribution to the target. This convergence is most often guaranteed through detailed balance, a sufficient condition for the chain to have the target equilibrium distribution. In practice, for any proposal distribution, one can ensure detailed balance through a Metropolis-Hastings (Hastings, 1970) accept/reject step.

Despite theoretical guarantees of eventual convergence, in practice convergence and mixing speed depend strongly on choosing a proposal that works well for the task at hand. What’s more, it is often more art than science to know when an MCMC chain has converged (“burned-in”), and when the chain has produced a new uncorrelated sample (“mixed”). Additionally, the reliance on detailed balance, which assigns equal probability to the forward and reverse transitions, often encourages random-walk behavior and thus slows exploration of the space (Ichiki & Ohzeki, 2013).

For densities over continuous spaces, Hamiltonian Monte Carlo (HMC; Duane et al. [1987] Neal, [2011]) introduces independent, auxiliary momentum variables, and computes a new state by integrating Hamiltonian dynamics. This method can traverse long distances in state space with a single Metropolis-Hastings test. This is the state-of-the-art method for sampling in many domains. However, HMC can perform poorly in a number of settings. While HMC mixes quickly spatially, it struggles at mixing across energy levels due to its volume-preserving dynamics. HMC also does not work well with multi-modal distributions, as the probability of sampling a large enough momentum to traverse a very low-density region is negligibly small. Furthermore, HMC struggles with ill-conditioned energy landscapes (Girolami & Calderhead, 2011) and deals poorly with rapidly changing gradients (Sohl-Dickstein et al., 2014).

Recently, probabilistic models parameterized by deep neural networks have achieved great success at modeling highly complex, multi-modal distributions (Kingma & Welling, 2013; Goodfellow et al., 2014). Building on these successes, we present a method that, given an analytically described distribution, automatically returns a sampler with good convergence and mixing properties, from a class of highly expressive parametric models. The proposed family of samplers is a generalization of HMC; it transforms the HMC trajectory using parametric functions (deep networks in our experiments), while conserving theoretical guarantees with a tractable Metropolis-Hastings accept/reject test.
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step. The sampler is trained to minimize a variation of expected squared jumped distance (similar in spirit to Pasarica & Gelman (2010)). Our parameterization reduces easily to standard HMC. It is further capable of emulating several common extensions of HMC such as within-trajectory tempering and diagonal mass matrices.

We evaluate our method on distributions where HMC usually struggles, as well as on a the real-world task of training latent-variable generative models.

Our contributions are as follows:

- We introduce a generic training procedure which takes as input a distribution defined by an energy function, and returns a fast-mixing MCMC kernel.
- We show significant empirical gains on various distributions where HMC performs poorly.
- We finally evaluate our method on the real-world task of training and sampling from a latent variable generative model, where we show improvement in the model’s log-likelihood, and greater complexity in the distribution of posterior samples.

2 Related Work

Adaptively modifying proposal distributions to improve convergence and mixing has been explored in the past (Andrieu & Thoms, 2008). In the case of HMC, prior work has reduced the need to choose step size (Neal, 2011) or number of leapfrog steps (Hoffman & Gelman, 2014) by adaptively tuning those parameters. Salimans et al. (2015) proposed an alternate scheme based on variational inference. We adopt the much simpler approach of Pasarica & Gelman (2010), who show that choosing the hyperparameters of a proposal distribution to maximize expected squared jumped distance is both principled and effective in practice. Most similar to ours is recent work from Song et al. (2017), which uses adversarial training of a volume-preserving transformation, which is subsequently used as an MCMC proposal distribution. While promising, this technique has several limitations. It does not use gradient information, which is often crucial to maintaining high acceptance rates, especially in high dimensions. It also can only indirectly measure the quality of the generated sample using adversarial training, which is notoriously unstable, suffers from “mode collapse” (where only a portion of a target distribution is covered), and often requires objective modification to train in practice (Arjovsky et al., 2017). Finally, since the proposal transformation preserves volume, it can suffer from the same difficulties in mixing across energy levels as HMC, as we illustrate in Section 5.

To compute the Metropolis-Hastings acceptance probability for a deterministic transition, the operator must be invertible and have a tractable Jacobian. Recent work (Dinh et al., 2016), introduces RNVP, an invertible transformation that operates by, at each layer, modifying only a subset of the variables by a function that depends solely on the remaining variables. This is exactly invertible with an efficiently computable Jacobian. Furthermore, by chaining enough of these layers, the model can be made arbitrarily expressive. This parameterization will directly motivate our extension of the leapfrog integrator in HMC.

3 Background

3.1 MCMC Methods and Metropolis-Hastings

Let $p$ be a target distribution, analytically known up to a constant, over a space $X$. MCMC methods (Neal, 1993) aim to provide samples from $p$. To that end, MCMC methods construct a Markov Chain whose stationary distribution is the target distribution $p$. Obtaining samples then corresponds to simulating a Markov Chain, i.e., given an initial distribution $\pi_0$ and a transition kernel $K$, constructing the following sequence of random variables:

$$
X_0 \sim \pi_0, \quad X_{t+1} \sim K(\cdot | X_t)
$$

In order for $p$ to be the stationary distribution of the chain, there are two conditions: $K$ must be irreducible and aperiodic and, more importantly, $p$ has to be a fixed point of $K$. The second condition can be expressed as: $p(x') = \int K(x' | x) p(x) dx$. This condition is most often satisfied by satisfying the stronger detailed balance condition, which can be written as: $p(x')K(x | x') = p(x)K(x' | x)$.
Given any proposal distribution \( q \), satisfying mild conditions, we can easily construct a transition kernel that respects detailed balance using Metropolis-Hastings \cite{Hastings1970} accept/reject rules. More formally, starting from \( x_0 \sim p_0 \), at each step \( t \), we sample \( x' \sim q(\cdot|x_t) \), and with probability \( A(x'|x_t) = \min \left(1, \frac{p(x')q(x|\cdot)}{p(x)q(x'|\cdot)} \right) \), accept \( x' \) as the next sample \( x_{t+1} \) in the chain. If we reject \( x' \), then we retain the previous state and \( x_{t+1} = x_t \). For typical proposals this algorithm has strong asymptotic guarantees. But in practice one must often choose between very low acceptance probabilities and very cautious proposals, both of which lead to slow mixing. For continuous state spaces, Hamiltonian Monte Carlo (HMC; \cite{Neal2011}) tackles this problem by proposing updates that move far in state space while staying roughly on iso-probability contours of \( p \).

### 3.2 Hamiltonian Monte Carlo

Without loss of generality, we assume \( p(x) \) to be defined by an energy function \( U(x) \), s.t. \( p(x) \propto \exp(-U(x)) \), and where the state \( x \in \mathbb{R}^n \). HMC extends the state space with an additional momentum vector \( v \in \mathbb{R}^n \), where \( v \) is distributed independently from \( x \), as \( p(v) \propto \exp(-\frac{1}{2} v^T v) \) (i.e., identity-covariance Gaussian). From an augmented state \( \xi = \{x,v\} \), HMC produces a proposed state \( \xi' = \{x',v'\} \) by approximately integrating Hamiltonian dynamics jointly on \( x \) and \( v \), with \( U(x) \) taken to be the potential energy, and \( \frac{1}{2} v^T v \) the kinetic. Since Hamiltonian dynamics conserve the total energy of a system, their approximate integration moves along approximate iso-probability contours of \( p(x,v) = p(x)p(v) \).

Hamiltonian dynamics are typically integrated using the leapfrog integrator \cite{Hairer2003}, which for a single time step consists of:

\[
v_{t+1}^\frac{1}{2} = v - \frac{\epsilon}{2} \partial_x U(x); \quad x' = x + \epsilon v_{t+1}^\frac{1}{2}; \quad v' = v - \frac{\epsilon}{2} \partial_x U(x').
\]  

Following \cite{Sohl-Dickstein2014}, we write the action of the leapfrog integrator in terms of an operator \( L \), \( L \xi = L\{x,v\} = \{x',v'\} \), and introduce a momentum flip operator \( F \), \( F\{x,v\} = \{x,-v\} \). It is important to note two properties of these operators. First, the transformation \( FLFL \) is an involution, i.e. \( FLFL\{x,v\} = FL\{x',-v'\} = \{x,v\} \). Second, the transformations from \( \{x,v\} \) to \( \{x,v_{t+1}^\frac{1}{2}\} \), from \( \{x,v_{t+1}^\frac{1}{2}\} \) to \( \{x',v_{t+1}^\frac{1}{2}\} \), and from \( \{x',v_{t+1}^\frac{1}{2}\} \) to \( \{x',v'\} \) are all volume-preserving shear transformations i.e., only one of the variables \( (x \text{ or } v) \) changes, by an amount determined by the other one. The determinant of the Jacobian, \( \left| \frac{\partial (FLFL)}{\partial \xi} \right| = 1 \), but we will leave it in symbolic form for use in Section 4. The Metropolis-Hastings acceptance probability for the HMC proposal is made simple by these two properties, and is

\[
A(FL\xi|\xi) = \min \left(1, \frac{p(FL\xi)p(\xi)}{p(\xi)p(FL\xi)} \left| \frac{\partial (FLFL)}{\partial \xi} \right| \right).
\]  

### 4 L2HMC: Training MCMC Samplers

In this section, we describe our proposed method L2HMC. Given access to only an energy function \( U \) (and not samples), L2HMC learns a parametric leapfrog operator \( L_\theta \) over an augmented state space. We begin by describing what desiderata we have for \( L_\theta \), then go into detail on how we parameterize our sampler. Finally, we conclude this section by describing our training procedure.

HMC is a powerful algorithm, but it can still struggle even on very simple problems. For example, a two-dimensional multivariate Gaussian with an ill-conditioned covariance matrix can take arbitrarily long to traverse (even if the covariance is diagonal), whereas it is trivial to sample directly from it. Another problem is that HMC cannot easily traverse low-density zones. For example, given a simple Gaussian mixture model, HMC cannot mix between modes without recourse to additional tricks, as illustrated in Figure \cite{Blei2014}. These observations determine the list of desiderata for our learned MCMC kernel: fast mixing, fast burn-in, mixing across energy levels, and mixing between modes.

While pursuing these goals, we must take care to ensure that our proposal operator retains two key features of the leapfrog operator used in HMC: it must be invertible, and the determinant of its Jacobian must be tractable. The leapfrog operator satisfies these properties by ensuring that each sub-update only affects a subset of the variables, and that no sub-update depends nonlinearly on any of the variables being updated. We are free to generalize the leapfrog operator in any way that
preserves these properties. In particular, we are free to translate and rescale each sub-update of the leapfrog operator, so long as we are careful to ensure that these translation and scale terms do not depend on the variables being updated.

As in HMC, we begin by augmenting the current state $x \in \mathbb{R}^n$ with a continuous momentum variable $v \in \mathbb{R}^n$ drawn from a standard normal. We also introduce a binary direction variable $d \in \{-1, 1\}$, drawn from a uniform distribution. Finally, to each step $t$ of the operator $L_0$ we assign a fixed random binary mask $m^t \in \{0, 1\}^n$ that will determine which variables are affected by each sub-update. We draw $m^t$ uniformly from the set of binary vectors satisfying $\sum_{i=1}^n m^t_i = \lfloor \frac{n}{2} \rfloor$, that is, half of the entries of $m^t$ are 0 and half are 1. For convenience, we write $\tilde{m}^t = 1 - m^t$ and $x_{m^t} = x \odot m^t$ ($\odot$ denotes element-wise multiplication). We will denote the complete augmented state as $\xi \triangleq \{x, v, d\}$, with probability density $p(\xi) = p(x)p(v)p(d)$.

We now describe the details of our augmented leapfrog integrator $L_0$, for a single time-step $t$, and for direction $d = 1$.

We first update the momenta $v$. This update can only depend on a subset $\zeta_1 \triangleq \{x, \partial_x U(x), t\}$ of the full state, which excludes $v$. It takes the form

$$v' = v \odot \exp\left(\frac{1}{2} S_v(\zeta_1)\right) - \frac{1}{2} \left(\partial_x U(x) \odot \exp(Q_v(\zeta_1)) + T_v(\zeta_1)\right).$$

(4)

We have introduced three new functions of $\zeta_1$: $T_v, Q_v, S_v$. $T_v$ is a translation, $\exp(Q_v)$ rescales the gradient, and $\exp\left(\frac{1}{2} S_v\right)$ rescales the momentum. The determinant of the Jacobian of this transformation is $\exp\left(\frac{1}{2} \cdot S_v(\zeta_1)\right)$. Note that if $T_v, Q_v,$ and $S_v$ are all zero, then we recover the standard leapfrog momentum update.

We now update $x$. As hinted above, to make our transformation more expressive, we first update a subset of the coordinates of $x$, followed by the complementary subset. The first update, which affects only $x_{m^t}$, depends on the state subset $\zeta_2 \triangleq \{x_{m^t}, v, t\}$. Conversely, with $x'$ defined below, the second update only affects $x_{\tilde{m}^t}$ and depends only on $\zeta_3 \triangleq \{x'_{m^t}, v, t\}$:

$$x' = x_{m^t} + m^t \odot \left[x \odot \exp(\epsilon S_x(\zeta_2)) + \epsilon v' \odot \exp(Q_x(\zeta_2)) + T_x(\zeta_2)\right]$$

$$x'' = x'_{m^t} + \tilde{m}^t \odot \left[x' \odot \exp(\epsilon S_x(\zeta_3)) + \epsilon v' \odot \exp(Q_x(\zeta_3)) + T_x(\zeta_3)\right].$$

(5)

Again, $T_x$ is a translation, $\exp(Q_x)$ rescales the effect of the momenta, and $\exp(\epsilon S_x)$ rescales the positions $x$, and we recover the original leapfrog position update if $T_x = Q_x = S_x = 0$. The determinant of the Jacobian of the first transformation is $\exp(\epsilon m^t \cdot S_x(\zeta_2))$, and the determinant of the Jacobian of the second transformation is $\exp(\epsilon \tilde{m}^t \cdot S_x(\zeta_3))$.

Finally, we update $v$ again, based on the subset $\zeta_4 \triangleq \{x'', \partial_x U(x''), t\}$:

$$v'' = v' \odot \exp\left(\frac{1}{2} S_v(\zeta_4)\right) - \frac{1}{2} \left(\partial_x U(x'') \odot \exp(Q_v(\zeta_4)) + T_v(\zeta_4)\right).$$

(6)

This update has exactly the same form as the momentum update in equation (4).

To give intuition into these terms, the scaling applied to the momentum can enable, among other things, acceleration in low-density zones, to facilitate mixing between modes. The scaling term applied to the gradient of the energy may allow better conditioning of the energy landscape (e.g., by learning a diagonal inertia tensor), or partial ignoring of the energy for rapidly changing gradients.

The corresponding integrator for $d = -1$ is given in Appendix A; it essentially just inverts the updates in equations (4) and (5). For all experiments, the functions $Q, S, T$ are implemented using multi-layer perceptrons, with shared weights. We encode the current time step in the MLP input.

Our leapfrog operator $L_0$ corresponds to running $M$ steps of this modified leapfrog, $L_0 \xi = L_0 \{x, v, d\} = \{x^{t\times M}, v^{t\times M}, d\}$, and our flip operator $F$ reverses the direction variable $d$, $F \xi = \{x, v, -d\}$. Written in terms of these modified operators, our proposal and acceptance probability are identical to those for standard HMC. Note, however, that this parameterization enables learning non-volume-preserving transformations, as the determinant of the Jacobian is a function of $S_x$ and $S_v$ that does not necessarily evaluate to 1. This quantity is derived in Appendix B.

For convenience, we denote by $R$ an operator that re-samples the momentum and direction. I.e., given $\xi = \{x, v, d\}$, $R \xi = \{x, v', d'\}$ where $v' \sim \mathcal{N}(0, I), d' \sim U\{-1, 1\}$. Sampling thus consists of alternating application of the FL and R, in the following two steps each of which is a Markov transition for $p$: 
1. \( \xi' = \text{FL}_\theta \xi \) with probability \( A(\text{FL}_\theta \xi | \xi) \) (Equation [3]), otherwise \( \xi' = \xi \).
2. \( \xi' = R \xi \)

This parameterization is effectively a generalization of standard HMC as it is non-volume preserving, with learnable parameters, and easily reduces to standard HMC for \( Q, S, T = 0 \).

4.1 Loss and Training Procedure

We need some criterion to tune the parameters \( \theta \) that control the functions \( Q, S, \) and \( T \). We choose a loss designed to reduce mixing time. Specifically, we aim to minimize lag-one autocorrelation.

This is equivalent to maximizing expected squared jumped distance \((\text{Pasarica & Gelman, 2010})).\) For \( \xi, \xi' \) in the extended state space, we define \( \delta(\xi', \xi) = \delta(\{x', v', \alpha', \{x, v, d\} \} = ||x - x'||^2_2 \).

Expected squared jumped distance is thus \( E_{\xi \sim p(\xi)} \left[ \delta(\text{FL}_\theta \xi, \xi) A(\text{FL}_\theta \xi | \xi) \right]. \) However, this loss need not encourage mixing across the entire state space. Indeed, maximizing this objective can lead to regions of state space where almost no mixing occurs, so long as the average squared distance traversed remains high. To optimize both for typical and worst case behavior, we include a reciprocal term in the loss,

\[
\ell_\lambda(\xi, \xi', A(\xi' | \xi)) = \frac{\lambda^2}{\delta(\xi', \xi) A(\xi' | \xi)} - \frac{\delta(\xi', \xi) A(\xi' | \xi)}{\lambda^2},
\]

where \( \lambda \) is a scale parameter, capturing the dimension of the problem. The second term encourages typical moves to be large, while the first term strongly penalizes the sampler if it is ever in a state where it cannot move far. We train our sampler by minimizing this loss over both the target distribution and initialization distribution. Formally, given an initial distribution \( p_0 \) over \( \chi \), we define \( q(\xi) = p_0(x)N(v', 0, 1)p(d) \), and minimize

\[
L(\theta) \triangleq E_{p(\xi)} \left[ \ell_\lambda(\xi, \text{FL}_\theta \xi, A(\text{FL}_\theta \xi | \xi)) \right] + E_{q(\xi)} \left[ \ell_\lambda(\xi, \text{FL}_\theta \xi, A(\text{FL}_\theta \xi | \xi)) \right].
\]

The first term of this loss encourages mixing as it considers our operator applied on elements of the distribution; the second term is meant to reward fast burn-in. Given this loss, we exactly describe our training procedure in Algorithm 1. It is important to note that each training iteration can be done with only one pass through the network and can be efficiently batched. We further emphasize that this training procedure can be applied to any learnable operator whose Jacobian’s determinant is tractable, making it a general framework for training MCMC proposals.

Algorithm 1 Training L2HMC

Input: Energy function \( U : \chi \rightarrow \mathbb{R} \) and its gradient \( \nabla_x U : \chi \rightarrow \chi \), initial distribution over the augmented state space \( q \), number of iterations \( n_{\text{iters}} \), number of leapfrogs \( M \), learning rate \( \alpha_t \), \( \alpha_i \) \( i \leq n_{\text{iters}} \), batch size \( N \) and scale parameter \( \lambda \).

Initialize the parameters of the sampler \( \theta \).

Initialize \( \{\xi_p^{(i)}\}_{i \leq N} \) from \( q(\xi) \).

for \( t = 0 \) to \( n_{\text{iters}} - 1 \) do

Sample a minibatch \( \{\xi_q^{(i)}\}_{i \leq N} \) from \( q(\xi) \).

\( L \leftarrow 0 \)

for \( i = 1 \) to \( N \) do

\( \xi_p^{(i)} \leftarrow R \xi_p^{(i)} \)

\( \xi_p^{(i)} \leftarrow L + \ell_\lambda \left( \xi_p^{(i)}, \text{FL}_\theta \xi_p^{(i)}, A(\text{FL}_\theta \xi_p^{(i)} | \xi_p^{(i)}) \right) + \ell_\lambda \left( \xi_q^{(i)}, \text{FL}_\theta \xi_q^{(i)}, A(\text{FL}_\theta \xi_q^{(i)} | \xi_q^{(i)}) \right) \)

end for

end for

5 Experiments

We present an empirical evaluation of our trained sampler on a diversity of energy functions. We first present results on a collection of toy distributions capturing common pathologies of energy
landscapes, followed by results on a task from machine learning: maximum-likelihood training of deep generative models. For each, we compare against a well-tuned HMC and show significant gains in mixing time.

5.1 Varied Collection of Energy Functions

We evaluate our L2HMC sampler on a diverse collection of energy functions, each providing different challenges for standard HMC.

Ill-Conditioned Gaussian (ICG): Gaussian distribution with diagonal covariance spaced log-linearly between $10^{-2}$ and $10^2$. This is representative of L2HMC learning a diagonal inertia tensor.

Strongly correlated Gaussian (SCG): We rotate a diagonal Gaussian with covariance $[10^2, 10^{-2}]$ by $\frac{\pi}{4}$. This is an extreme version of an example from Neal (2011). This problem shows that, although our parametric sampler only applies element-wise transformations, it can learn rotations.

Mixture of Gaussians (MoG): Mixture of two isotropic Gaussians with $\sigma^2 = 0.1$, and centroids separated by distance 4. The means are thus about 12 standard deviations apart, making it almost impossible for HMC to mix between modes. This experiment shows that L2HMC can learn to do within-trajectory tempering (Neal, 2011).

Rough Well: Similar to an example from Sohl-Dickstein et al. (2014), for a given $\eta > 0$, $U(x) = \frac{1}{2} x^T x + \eta \sum_i \cos(\frac{x_i}{\eta})$. While the energy itself is altered negligibly, its gradient is perturbed by a high frequency noise oscillating between $-1$ and $1$. In our experiments, we choose $\eta = 10^{-2}$.

For each of these distributions, we compare against HMC with the same number of leapfrog steps and a well-tuned step-size. To compare mixing time, we plot auto-correlation for each method and report effective sample size (ESS), computed as described by Hoffman & Gelman (2014). We observe that samplers trained with L2HMC show greatly improved autocorrelation and ESS on the right mode to a smaller volume in the right mode. It is important to note that the distance between both clusters minus their standard deviations is less than in the left mode to a smaller volume in the left mode, making it almost impossible to traverse. This is also an issue for HMC. On the other hand, L2HMC can both traverse the low-density region between modes, and map a larger volume in the left mode to a smaller volume in the right mode. It is important to note that the distance between both clusters minus their standard deviations is less than in the MoG case, and it is thus a good diagnostic of the shortcomings of volume-preserving transformations.

5.2 Latent-Variable Generative Model

We apply our learned sampler to the task of training, and sampling from the posterior of, a latent-variable generative model. The model consists of a latent variable $z \sim p(z)$, where we choose $p(z) = \mathcal{N}(z; 0, I)$, and a conditional distribution $p(x|z)$ which generates the image $x$. Given a family of parametric ‘decoders’ $\{z \mapsto p(x|z; \phi), \phi \in \Phi\}$, and a set of samples $D = \{x^{(i)}\}_{i \in \mathcal{N}}$, training involves finding $\phi^* = \arg \max_{\phi \in \Phi} p(D; \phi)$. However, the log-likelihood is intractable as $p(x; \phi) = \int p(x|z; \phi)p(z)dz$. To remedy that problem, Kingma & Welling (2013) proposed jointly training an approximate posterior $q_\psi$ that maximizes a tractable lower-bound on the log-likelihood:

$$\mathcal{L}_{\text{ELBO}}(x, \phi, \psi) = \mathbb{E}_{q_\psi(z|x)}[p(x|z; \phi)] - \text{KL}(q_\psi(z|x)||p(z)) \leq p(x) \quad (9)$$

Recently, to improve upon well-known pitfalls like over-pruning (Yeung et al., 2017) of the VAE, Hoffman (2017) proposed HMC-DLGM. For a data sample $x^{(i)}$, after obtaining a sample from the approximate posterior $q_\psi(\cdot|x^{(i)})$, Hoffman (2017) runs a MCMC algorithm with energy function
To show the benefits of L2HMC, we borrow the method from Hoffman (2017), but replace HMC by L2HMC.

In the subsequent sections, we compare our method to the standard VAE model from Kingma & Welling (2013) and HMC-DGLM from Hoffman (2017). It is important to note that, since our sampler enables a more expressive, non-Gaussian, posterior.

Implementation details: Our decoder ($p_\phi$) is a neural network with 2 fully connected layers, with 1024 units each and softplus non-linearities, and outputs Bernoulli activation probabilities for each pixel. The encoder ($q_\psi$) has the same architecture, returning mean and variance for the approximate posterior. Our model was trained for 300 epochs with Adam (Kingma & Ba 2014) and a learning rate $\alpha = 10^{-3}$. All experiments were done on the dynamically binarized MNIST dataset (LeCun).

5.2.1 Sample Quality and Data Likelihood

We first present samples from decoders trained with L2HMC, HMC and the ELBO (i.e. vanilla VAE). Although higher log likelihood does not necessarily correspond to better samples (Theis et al. 2015), we can see in Figure 5 shown in the Appendix, that the decoder trained with L2HMC generates sharper samples than the compared methods.
We now compare our method to HMC in terms of log-likelihood of the data. As we previously stated, the marginal likelihood of a data point $x \in \mathcal{X}$ is not tractable as it requires integrating $p(x, z)$ over a high-dimensional space. However, we can estimate it using annealed importance sampling (AIS; Neal (2001)). Following Wu et al. (2016), we evaluate our generative models on both training and held-out data. In Figure 2, we plot the data’s log-likelihood against the number of gradient computation steps for both HMC-DGLM and L2HMC-DGLM. We can see that for a similar number of gradient computations, L2HMC-DGLM achieves higher likelihood for both training and held-out data. This is a strong indication that L2HMC provides significantly better posterior samples.

5.2.2 Increased Expressivity of the Posterior

In the standard VAE framework, approximate posteriors are often parametrized by a Gaussian, thus making a strong assumption of uni-modality. In this section, we show that using L2HMC to sample from the posterior enables learning of a richer posterior landscape.

**Block Gibbs Sampling** To highlight our ability to capture more expressive posteriors, we in-paint the top of an image using Block Gibbs Sampling using the approximate posterior or L2HMC. Formally, let $x_0$ be the starting image. We denote top or bottom-half pixels as $x_{\text{top}}^0$ and $x_{\text{bottom}}^0$. At each step $t$, we sample $z^{(t)} \sim p(z|x_t; \theta)$, sample $\tilde{x} \sim p(x|z_t; \theta)$. We then set $x_{\text{top}}^{t+1} = \tilde{x}_{\text{top}}$ and $x_{\text{bottom}}^{t+1} = x_{\text{bottom}}^t$. We compare the results obtained by sampling from $p(z|x; \theta)$ using $q_\psi$ (i.e. the approximate posterior) vs. our trained sampler. The results are reported in Figure 3a. We can see that L2HMC easily mixes between modes (3, 5, 8, and plausibly 9 in the figure) while the approximate posterior gets stuck on the same reconstructed digit (8 in the figure).

**Visualization of the posterior** After training a decoder with L2HMC, we randomly choose an element $x_0 \in D$ and run 512 parallel L2HMC chains for 20, 000 Metropolis-Hastings step. We then find the direction of highest variance, project the samples along that direction and show a histogram in Figure 3b. This plot shows non-Gaussianity in the latent space for the posterior. Using our improved sampler thus effectively enables the decoder to have a more expressive posterior.

6 Conclusion

In this work, we presented a general method to train expressive MCMC kernels parameterized with deep neural networks. Given a target distribution $p$, analytically known up to a constant, our method provides a fast-mixing sampler, able to efficiently explore the state space. Our hope is that our method can be utilized in a “black-box” manner, in domains where sampling constitutes a huge bottleneck such as protein foldings (Schütte et al., 1999) or physics simulations (Olsson, 1995). To this effect, we will open-source our code to the community.
REFERENCES


Appendix

A REVERSE LEAPFROG OPERATOR

Let’s $\xi = \{x, v, d\}$ in the extended state space with $d = -1$. Here, we describe the leapfrog updates for a single time step $t$, this consists of inverting the equations presented in the corresponding section.

Let $\zeta_1 = \{x, v, t\}$, we have:

$$v' = \left[ x + \frac{\epsilon}{2} \left( \partial_x U(x) \circ \exp(Q_v(\zeta_1)) + T_v(\zeta_1) \right) \right] \circ \exp(-S_v(\zeta_1)).$$

With the notation from Section 4, let $\zeta_2 \triangleq \{x_{mt}, v, t\}$

$$x' = x_{mt} + \hat{m} \circ \left[ (x - \epsilon \left( \exp(Q_x(\zeta_2)) \circ v' + T_x(\zeta_2) \right)) \circ \exp(-\epsilon S_v(\zeta_2)) \right].$$

Let us denote $\zeta_3 \triangleq \{x', m, t\}$:

$$x' = x_{mt} + m \circ \left[ (x' - \epsilon \left( \exp(Q_x(\zeta_3)) \circ v' + T_x(\zeta_3) \right)) \circ \exp(-\epsilon S_v(\zeta_3)) \right].$$

Finally, the last update, with $\zeta_4 \triangleq \{x'', \partial_x U(x''), t\}$:

$$v' = \left[ x + \frac{\epsilon}{2} \left( \partial_x U(x'') \circ \exp(Q_v(\zeta_4)) + T_v(\zeta_4) \right) \right] \circ \exp(-S_v(\zeta_4)).$$

It is important to note that to invert $L_\theta$, these steps should be ran for $t$ from $M$ to 1.

B DETERMINANT OF THE JACOBIAN

Given the derivations (and notations) from Section 4 for the forward operator $L_\theta$, we can immediately compute the Jacobian:

$$\log \left| \partial[FL_\theta(\xi)] \frac{\partial}{\partial \xi^T} \right| = \sum_{t \leq M} \left[ \frac{\epsilon}{2} \cdot S_v(\zeta_1^t) + \epsilon m^t \cdot S_x(\zeta_2^t) + \epsilon \hat{m}^t \cdot S_x(\zeta_3^t) + \frac{\epsilon}{2} \cdot S_v(\zeta_4^t) \right].$$

Where $\zeta_i^t$ denotes the intermediary variable $\zeta_i$ at time step $t$ and $d$ is the direction of $\xi$ i.e. $\xi = \{x, v, d\}$.

C IMPLEMENTATION DETAILS OF L2HMC

First of all, we keep separate parameters for the network responsible for updating $v$ and those updating $x$. The architectures are the same. Let us take the example of $Q_v, S_v, T_v$. The time step $t$ is given as input to the MLP, encoded as $\tau(t) = (\cos(\frac{2\pi t}{M}), \sin(\frac{2\pi t}{M}))$. $\sigma(\cdot)$ denotes the ReLU non-linearity.

For $n_h$ hidden units per layer:

- We first compute $h_1 = \sigma(W_1 x + W_2 v + W_3 \tau(t) + b) \in \mathbb{R}^{n_h}$.
- $h_2 = \sigma(W_h h_1 + b_1) \in \mathbb{R}^{n_h}$.
- $S_v = \lambda_v \tanh(W_v h_2 + b_v), Q_v = \lambda_v \tanh(W_q h_2 + b_q), T_v = W_v h_2 + b_t$.

In Section 5.1 the $Q, S, T$ are neural networks with 2 hidden layers with 10 (100 for the 50-d ICG) units and ReLU non-linearities. We train with Adam [Kingma & Ba 2014] and a learning rate $\alpha = 10^{-3}$. We train for 5,000 iterations with a batch size of 200.

For the MoG tasks, we train our sampler with a temperature parameter that we continuously anneal; we evaluate the trained sampler without using temperature.
In addition, the auxiliary variable $x$ (here the image from MNIST) is first passed through a 2-layer neural network, with softplus non-linearities and 512 hidden units. This input is given to both networks $\{\cdot\}_x$ and $\{\cdot\}_v$. The architecture then consists of 2 hidden layers of 200 units and ReLU non-linearities. For $\lambda$ (scale parameter of the loss), we use the standard deviation of the approximate posterior.

D.2 IMPLEMENTATION DETAILS OF L2HMC-DGLM

Similar to our L2HMC training on unconditional sampling, we share weights across $Q, S$ and $T$. In addition, the auxiliary variable $x$ is first passed through a 2-layer neural network, with softplus non-linearities and 512 hidden units. This input is given to both networks $\{\cdot\}_x$ and $\{\cdot\}_v$. The architecture then consists of 2 hidden layers of 200 units and ReLU non-linearities. For $\lambda$ (scale parameter of the loss), we use the standard deviation of the approximate posterior.

D.1 TRAINING ALGORITHM

In this section is presented our training algorithm as well as a diagram explaining L2HMC-DGLM. For conciseness, given our operator $L_\theta$, we denote by $K_\theta(\cdot|x)$ the distribution over next state given sampling of a momentum and direction and the Metropolis-Hastings step.

Algorithm 2 L2HMC for latent variable generative models

**Input:** dataset $D$, number of iterations $n_{\text{iters}}$, number of Metropolis-Hastings step $J$, number of leapfrogs $M$, and learning rate schedule $(\alpha_t)_{t \leq n_{\text{iters}}}$.

Randomly initialize the decoder’s parameters $\phi$ and the approximate posterior $\psi$. Initialize the parameters of the sampler $\theta$ with $M$ leapfrog steps.

for $t = 0$ to $n_{\text{iters}} - 1$

Randomly sample a minibatch $B$ from the dataset $D$.

$L_{\text{ELBO}}, L_{\text{Sampler}}, L_{\text{Decoder}} \leftarrow 0$

for $x^{(b)} \in B$ do

Sample $\xi_0^{(b)} \sim r(.|x^{(b)}; \psi)$.

$L_{\text{ELBO}} \leftarrow p(x^{(b)}|z_0^{(b)}; \phi) - \text{KL}(q_\psi(z|x^{(b)}))$ with $\xi_0^{(b)} = \{z_0^{(b)}, v_0^{(b)}, d_0^{(b)}\}$

Define the energy function $U_x(z) = -\log p(x|z; \theta) - \log p(z)$

$L_{\text{Sampler}} \leftarrow 0$

$\lambda \leftarrow \sqrt{\text{Var}(q_\psi(z_0^{(b)}|x^{(b)}))}$

for $j = 0$ to $J - 1$ do

$\xi_j^{(b)} \leftarrow R\xi_j^{(b)}$

$L_{\text{Sampler}} \leftarrow L_{\text{Sampler}} + \ell_\lambda(\xi_j^{(b)}, \text{FL}_\theta\xi_j^{(b)}, \text{A}(\text{FL}_\theta\xi_j^{(b)}|\xi_j^{(b)}))$

Set $\xi_j^{(b)}$ to $\text{FL}_\theta\xi_j^{(b)}$ with probability $\text{A}(\text{FL}_\theta\xi_j^{(b)}|\xi_j^{(b)})$.

end for

$L_{\text{Decoder}} \leftarrow L_{\text{Decoder}} + \log p(x^{(b)}|z_j^{(b)}; \phi)$ with $\xi_j^{(b)} = \{z_j^{(b)}, v_j^{(b)}, d_j^{(b)}\}$

end for

$\phi \leftarrow \phi + \alpha_t \nabla_\phi L_{\text{Decoder}}$

$\psi \leftarrow \psi + \alpha_t \nabla_\psi L_{\text{ELBO}}$

$\theta \leftarrow \theta + \alpha_t \nabla_\theta L_{\text{Sampler}}$

end for

Figure 4: Diagram of our L2HMC-DGLM model. Nodes are functions of their parents. Round nodes are deterministic, diamond nodes are stochastic and the doubly-circled node is observed.
AIS Evaluation  For each data point, we run 20 Markov Chains in parallel, 10,000 annealing steps with 10 leapfrogs per step and choose the step size for an acceptance rate of 0.65.

D.3 MNIST Samples