Towards a Complete Analysis of Langevin Monte Carlo: Beyond Poincaré Inequality

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Abstract

Langevin diffusions are rapidly convergent under appropriate functional inequality assumptions. Hence, it is natural to expect that with additional smoothness conditions to handle the discretization errors, their discretizations like the Langevin Monte Carlo (LMC) converge in a similar fashion. This research program was initiated by [VW19], who established results under log-Sobolev inequalities. $[CEL^+22]$ extended the results to handle the case of Poincaré inequalities. In this paper, we go beyond Poincaré inequalities, and push this research program to its limit. We do so by establishing upper and lower bounds for Langevin diffusions and LMC under weak Poincaré inequalities that are satisfied by a large class of densities including polynomially-decaying heavy-tailed densities (i.e., Cauchy-type). Our results explicitly quantify the effect of the initializer on the performance of the LMC algorithm. In particular, we show that as the tail goes from sub-Gaussian, to sub-exponential, and finally to Cauchy-like, the dependency on the initial error goes from being logarithmic, to polynomial, and then finally to being exponential. This three-step phase transition is in particular unavoidable as demonstrated by our lower bounds, clearly defining the boundaries of LMC.

1 Introduction

Consider the problem of sampling from a target probability density $\pi \propto \exp(-V)$ on \mathbb{R}^d using the canonical algorithm, Langevin Monte Carlo (LMC). The LMC iterations are given by

$$x_{k+1} = x_k - h\nabla V(x_k) + \sqrt{2h\xi_k},\tag{LMC}$$

where h > 0 is the step size, and $(\xi_k)_{k \in \mathbb{N}}$ is an i.i.d. sequence of standard Gaussian random vectors. This algorithm is based on discretizing the following stochastic differential equation (SDE), often referred to as the (overdamped) Langevin diffusion,

$$dX_t = -\nabla V(X_t)dt + \sqrt{2} \, dB_t,\tag{LD}$$

where $(B_t)_{t \in \mathbb{R}_+}$ is the *d*-dimensional standard Brownian motion. When π is (strongly) log-concave and smooth, non-asymptotic convergence of LMC has been extensively studied [DT12, Dal17b, Dal17a, DM19, DMM19].

The Langevin diffusion (LD), however, converges under relatively milder functional inequality assumptions which are less restrictive compared to global curvature conditions like log-concavity. Indeed, while log-concavity restricts π to be uni-modal, functional inequality based conditions allow for some degree of multi-modality in

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 π [CCNW21]. Furthermore, functional inequalities characterize a wide range of target densities by capturing the tail behavior of the potential. For example, a target potential with tail growth $V(x) \approx ||x||^{\alpha}$ at infinity, would satisfy a logarithmic Sobolev inequality (LSI) when $\alpha = 2$, and satisfies a Poincaré inequality (PI) when $\alpha = 1$. Thus, an LSI induces a faster tail growth and is consequently a stronger condition than a PI.

Motivated by this, the following research program was initiated by [VW19]: Can one provide convergence guarantees for (LMC) when the target density satisfies a functional inequality and a smoothness condition? The authors answered the question in the affirmative, showing that the following two conditions on the target $\pi \propto e^{-V}$ suffice to establish a sharp non-asymptotic guarantee for LMC: (i) π satisfies an LSI and (ii) ∇V is Lipschitz continuous. [CEL⁺22] extended this framework significantly; among other contributions, they also proved that LSI can be replaced with a Latała-Oleszkiewicz inequality (LOI), which can cover a range of tail behavior, i.e. $\alpha \in [1, 2]$, interpolating between the edge cases PI ($\alpha = 1$) and LSI ($\alpha = 2$). These works provide a thorough characterization of the convergence of LMC for at least linearly growing potentials, and to our knowledge, providing the state of the art guarantees under minimal set of conditions for this algorithm. However, it is rather unclear how much further this program can be extended. For example, what is the threshold for the tail behavior α beyond which LMC fails, if at all it fails to sample from such heavy-tailed targets?

In this paper, we aim to complete the program initiated by [VW19], and push the convergence analysis of LMC to its limits. We study the behavior of LMC for potentials that satisfy a family of weak-Poincaré inequalities (WPI), which are one of the mildest conditions required to prove the ergodicity of the Langevin diffusion [RW01, BGL14]. A particularly interesting aspect of WPI is that virtually any target density satisfies such an inequality. Thus, by proving a convergence guarantee for LMC under a WPI with explicit rate estimates, we establish its convergence universally for any sufficiently smooth target. Interestingly, for targets with sublinear tails, i.e. $V(x) \approx ||x||^{\alpha}$ for $\alpha \in (0, 1)$, our rate is polynomial in the initial error, and smoothly extrapolates the rate derived by [CEL⁺22] which was originally covering the regime $\alpha \in [1, 2]$. In the case α approaches 0, however, the tail is logarithmic $V(x) \approx \ln(||x||)$, e.g. for Cauchy-type distributions, and our rate estimates exhibit an exponential dependence on the initial error; thus, when there is no warm-start available, LMC would require exponentially many iterations in the initial error for such targets with extreme heavy tails. We also provide a lower bound for LMC under a general tail-growth condition, proving that for Cauchy-type distributions, there is an initialization such that this exponential dependence is unavoidable. Our main contributions can be summarized as follows.

- For a target $\pi \propto e^{-V}$ satisfying a WPI with a Hölder continuous ∇V , we establish non-asymptotic convergence guarantees for LMC and the Langevin diffusion in Rényi divergence. Since any distribution with a locally bounded potential V satisfies a WPI [RW01], our results provide a convergence guarantees for LMC for any sufficiently smooth target.
- We prove WPIs with explicit dimension dependence for two model examples of heavy-tailed distributions that do not satisfy a Poincaré inequality, hence cannot be covered by the results of $[\text{CEL}^+22]$. First, we consider sub-linearly decaying potentials of the form $V(x) = (1 + ||x||^2)^{\alpha/2}$ and establish a rate, which coincides with the estimates of $[\text{CEL}^+22]$, but also holds for $\alpha \in (0, 1)$, namely beyond a Poincaré inequality. Notably, this rate is polynomial in the initial error for all $\alpha > 0$. We also consider the case of extreme heavy tails, i.e. Cauchy-type potentials with $\nu > 0$ degrees of freedom of the form $V(x) = \frac{d+\nu}{2} \ln(1+||x||^2)$, which does not have moments of order $\geq \nu$ defined. For this class of distributions, we prove that, even though LMC converges in Rényi divergence of any order, the dependence on the initial error may be exponential, which may limit its performance severely.
- Finally, we establish lower bounds for the complexity of LMC as well as the Langevin diffusion in Rényi divergence, under various tail growths in the range $\alpha \in [0, 2]$. Our lower bounds indicate that, as the tail growth becomes heavier, LMC and the diffusion both exhibit a slow start behavior by having a worse dependence on the initial divergence. In the particular case of Cauchy-type targets, the exponential dependence on the initial error for LMC and the diffusion is unavoidable, unless there is a good initialization available.

More related work. There have been innumerable works in the recent past focusing on (strongly) log-concave sampling with LMC, which makes it hard to summarize them here. We refer the interested reader to [Che23] for a detailed exposition. Beyond the log-concave setting, the assumption of dissipativity, which controls the growth order of the potential, is used in a large number of prior works to obtain convergence rates for LMC [DM17, RRT17, EMS18, EH21, MFWB22, EHZ22]. Additionally, a recent result by [BCE⁺22] developed a first-order stationarity theory for (averaged) LMC by providing Fisher information gaurantess for sampling from target densities that are only Hölder continuous without any further functional inequality or curvature-based assumptions. When combined with standard transportation inequalities, their results provide convergence guarantees in the TV distance for the averaged LMC under under the assumption that the target satisfies Poincaré inequality.

We also remark that in the (strongly) log-concave or light-tailed settings, several non-asymptotic results exist on variants of LMC, including higher order integrators [SL19, LWME19, HBE20], the underdamped Langevin Monte Carlo [CCBJ18, EGZ19, CLW19, DRD20, ZCL⁺23], the Metropolis-adjusted Langevin Algorithm [DCWY18, LST20, CLA⁺21, WSC22], and the Riemannian Langevin Algorithm [LE20].

Research on the analysis of heavy-tailed sampling is relatively scarce, especially results that are nonasymptotic in nature. [CDV09] studied the iteration complexity of the Metropolis random walk algorithm for sampling from s-concave distributions. [JR07] established polynomial ergodicity results for several sampling algorithms including LMC. [Kam18] developed modifications of the standard Metropolis random walk that are suitable for handling heavy-tailed targets and established asymptotic convergence results. Recently, [ALPW22] analyzed Metropolis random walk algorithms under WPI and established rate of convergence results in variance-like metrics. [JG12] introduced a variable transformation method for Metropolis random walk algorithms, transforming heavy-tailed densities into light-tailed ones using invertible transformations to benefit from existing light-tailed sampling algorithms. [HBE22] looked into ULA on a class of transformed densities and provided non-asymptotic results, mainly focusing on isotropic densities. The transformation approach has been extended in recent works such as [YLR22], and has also been used to prove asymptotic exponential ergodicity for various sampling algorithms in the heavy-tailed settings [DBCD19, DGM20, BRZ19]. [LWME19] and [HFBE23] analyze the iteration complexity of sampling by discretizing certain Itô diffusions.

While the literature on upper bounds on the complexity of sampling algorithms has seen significant progress, the literature on lower bounds is quite limited. Algorithm-independent query complexity of sampling from strongly log-concave distributions in one dimension was obtained by [CGL⁺22]. [LZT22] established lower bounds for LMC for sampling from strongly log-concave distributions. [CBL22] established lower bounds for sampling from strongly log-concave distributions in the stochastic setting, when the gradients are observed with noise. [GLL20] established lower bounds for the related problem of estimating the normalizing constants of a log-concave density. [LST21] and [WSC22] established lower bounds for the class of metropolized algorithms (including metropolized Langevin and Hamiltonian Monte Carlo methods) for sampling from strongly log-concave distributions. Finally, lower bounds in Fisher information for non-log-concave sampling were obtained in [CGLL22].

Notation. Throughout the paper, we will use $\pi \propto \exp(-V)$ to denote the target probability measure with unnormalized potential V, and ρ_t and μ_k to denote the law of the Langevin diffusion at time t, and the law of the LMC at iteration k. $(P_t)_{t \in \mathbb{R}_+}$ will denote the Markov semigroup of the Langevin diffusion. Probability measures we work with in this paper admit densities with respect to the Lebesgue measure, and we will use the same notation for their densities. We will use $\tilde{\Theta}_{\Psi}$ and \tilde{O}_{Ψ} to hide polylog factors and constants depending only on the set of variables Ψ . $\Gamma(z) \coloneqq \int_0^\infty t^{z-1} e^{-t} dt$ for z > 0 denotes the Gamma function, and ω_d is the volume of the unit *d*-ball.

2 Weak Poincaré Inequalities and Rényi Convergence of the Diffusion

We consider the following class of functional inequalities introduced by [RW01], motivated by the work of $[Lig91]^1$.

Definition 1 (Weak Poincaré Inequality). A probability measure π on \mathbb{R}^d satisfies a weak Poincaré inequality (WPI) if there exists $\beta_{\mathsf{WPI}} : (0, \infty) \to \mathbb{R}_+$ and $\Phi : L^2(\pi) \to [0, \infty]$ with $\Phi(cf + a) = c^2 \Phi(f)$ for every $c, a \in \mathbb{R}$

 $^{^{1}}$ We refer the interested reader to [Aid98] and [Mat98] for other attempts to propose weaker versions of Poincaré inequalities, and to [RW01] for their relationship with Definition 1.

and $f \in L^2(\pi)$, such that for every smooth f,

$$\operatorname{Var}_{\pi}(f) \leq \beta_{\mathsf{WPI}}(r) \, \mathbb{E}_{\pi} \left[\left\| \nabla f \right\|^{2} \right] + r \, \Phi(f), \quad \forall r > 0.$$
(WPI)

We remark that virtually any target measure of interest in sampling satisfies such an inequality. More specifically, [RW01] showed that $\pi \propto e^{-V}$ satisfies a WPI with $\Phi(\cdot) = \operatorname{Osc}(\cdot)^2 := (\sup f - \inf f)^2$ and for some β_{WPI} as soon as V is locally bounded. In the special cases where $\Phi = 0$ or the function β_{WPI} is uniformly bounded, the above inequality reduces to the classical Poincaré inequality which reads, for a constant β_{PI} and for every smooth f,

$$\operatorname{Var}_{\pi}(f) \le \beta_{\mathsf{PI}} \mathbb{E}_{\pi} \left[\left\| \nabla f \right\|^{2} \right].$$
(PI)

The tail properties of the distribution π is captured by the function β_{WPI} , which will essentially determine the convergence rate of LMC. We will present our convergence guarantees under the generic condition (WPI), and for several model examples, we will derive explicit estimates of β_{WPI} to make our results more explicit.

Functional inequalities of the form (WPI) naturally extend Poincaré inequalities (PI) to arbitrary distributions, removing any tail growth requirement in a PI. In particular, as a PI is equivalent to an exponential L^2 convergence rate for the Markov semigroup, a WPI is equivalent to a subexponential L^2 convergence rate [RW01, BGL14]. Similarly, one can also replace the variance term in (WPI) with entropy, in which case the functional inequality is of the form of a weak log-Sobolev inequality (WLSI). As shown by [CGG07], a WLSI is equivalent to a WPI; thus, we find it sufficient to present our results in terms of the WPI.

Following the recent works (see, e.g., [GT20, EHZ22, CEL⁺22]), we use Rényi divergence as a measure of distance between two probability distributions. Rényi divergence of order q is defined by

$$R_q(\rho \parallel \pi) \coloneqq \frac{1}{q-1} \ln \left\| \frac{\mathrm{d}\rho}{\mathrm{d}\pi} \right\|_{L^q(\pi)}^q \quad \text{for} \quad 1 < q < \infty,$$

$$(2.1)$$

when ρ is absolutely continuous with respect to π , and $+\infty$ otherwise. By Jensen's inequality, $R_q(\rho \| \pi)$ is nondecreasing in q. If we consider the limits, (i) as $q \downarrow 1$ it reduces to KL divergence, i.e. $\lim_{q \downarrow 1} R_q(\rho \| \pi) = \mathsf{KL}(\rho \| \pi)$ and (ii) as $q \to \infty$ it reduces to the L^{∞} -norm, i.e. $\lim_{q \to \infty} R_q(\rho \| \pi) = \|\mathrm{d}\rho/\mathrm{d}\pi\|_{L^{\infty}(\pi)}$. It is also related to χ_2 divergence through $\chi_2(\rho \| \pi) + 1 = \exp(R_2(\rho \| \pi))$.

Providing convergence guarantees in Rényi divergence is of particular interest since it upper bounds many commonly used distance measures. Specifically, by Pinsker's inequality and the monotonicity of Rényi divergence, we have

$$2D_{\mathsf{TV}}(\rho,\pi)^2 \le \mathsf{KL}(\rho \parallel \pi) \le R_q(\rho \parallel \pi) \quad \text{for} \quad q > 1.$$

Notice that comparing the quadratic Wasserstein distance $W_2^2(\rho, \pi)$ with $R_q(\rho \parallel \pi)$ is more subtle. Under a PI with constant β_{PI} , and more broadly under finite fourth moments, one can write

$$\ln\left(1 + \frac{1}{2\beta_{\mathsf{Pl}}} \operatorname{W}_{2}^{2}(\rho, \pi)\right) \leq R_{2}(\rho \parallel \pi) \quad \text{and} \quad \ln\left(1 + \frac{\operatorname{W}_{2}^{4}(\rho, \pi)}{4\mathbb{E}_{\pi}\left[\|x\|^{4}\right]}\right) \leq R_{2}(\rho \parallel \pi),$$

respectively. The first inequality is due to [Liu20] and the second one can be derived from the weighted total variation control on W₂ [Vil03, Proposition 7.10] along with the Cauchy-Schwartz inequality. Hence, even under a WPI, a bound in Rényi divergence can be translated to a bound in W₂ (when $\mathbb{E}_{\pi}[||x||^4] < \infty$), possibly at the expense of introducing additional dimension dependency into the bounds. It is worth highlighting that a distribution satisfying (WPI) does not need to have any particular moment defined, in which case W₂ may be undefined, but its Rényi divergence of some order may still be well-defined.

2.1 Rényi Convergence of the Langevin Diffusion

Classically, convergence of the Langevin diffusion under (WPI) is known only under the variance metric, or equivalently, the χ^2 divergence (see, for *e.g.*, [Wan06, Chapter 4] and [BGL14, Chapter 7.5]). The following result characterizes its convergence in Rényi divergence which is stronger in the case of q > 2.

Theorem 2. Suppose π satisfies (WPI) for some β_{WPI} and $\Phi(\cdot) = \operatorname{Osc}(\cdot)^2$. For any $1 < q < q' \leq \infty$ such that $R_{q'}(\rho_0 \parallel \pi) < \infty$, define $\delta_0 := \exp(qR_{q'}(\rho_0 \parallel \pi))$. Then, for any r > 0,

$$R_{q}(\rho_{t} \parallel \pi) \leq \begin{cases} R_{q}(\rho_{0} \parallel \pi) - \frac{2 - 4r\delta_{0}}{\beta(r)q}t & \text{if } R_{q}(\rho_{0} \parallel \pi), R_{q}(\rho_{t} \parallel \pi) \geq 1\\ e^{-\frac{2t}{\beta(r)q}}(R_{q}(\rho_{0} \parallel \pi) - 2r\delta_{0}) + 2r\delta_{0} & \text{if } R_{q}(\rho_{0} \parallel \pi) < 1, \end{cases}$$

where

$$\beta(r) \coloneqq \begin{cases} \beta_{\mathsf{WPI}}(r) & \text{if } q' = \infty \\ \beta_{\mathsf{WPI}}\left((r/5)^{\frac{q'}{q'-q}}\right) \ln\left((5/r)^{\frac{q'}{q'-q}}\right) & \text{if } q' < \infty. \end{cases}$$
(2.2)

Therefore, we have $R_q(\rho_T \parallel \pi) \leq \varepsilon$ whenever

$$T \ge q\beta \left(\frac{1}{4\delta_0}\right) R_q(\rho_0 \parallel \pi) + \frac{q}{2}\beta \left(\frac{\varepsilon}{4\delta_0}\right) \ln\left(\frac{1}{\varepsilon}\right).$$
(2.3)

We emphasize that while the classical convergence results under (WPI) [Wan06, BGL14] require $R_{\infty}(\rho_0 \parallel \pi) < \infty$ at initialization, whereas our convergence guarantees hold as soon as the initial error satisfies $R_{q'}(\rho_0 \parallel \pi) < \infty$ for some q' > q. Moreover, in the case where π satisfies a PI, i.e. when β_{WPI} is constant, we can remove the requirement of $R_{q'}(\rho_0 \parallel \pi) < \infty$, and the above theorem recovers [VW19, Theorem 3] by defining $\beta(0) := \lim_{r \to 0} \beta_{WPI}(r)$.

The proof of Theorem 2 is presented in Appendix A.1, and it relies on a proof technique also used in [VW19, CEL⁺22]. However, because of the $\Phi(\cdot) = \operatorname{Osc}(\cdot)^2$ term in (WPI), we additionally need to control oscillations of $\frac{d\rho_t}{d\pi}$ uniformly over the process. Via contraction properties for the semigroup, we can reduce such control to $\frac{d\rho_0}{d\pi} \in L^{\infty}(\pi)$ (or equivalently $R_{\infty}(\rho_0 \parallel \pi) < \infty$). To further relax this assumption, one needs to obtain a WPI with a *weaker* Φ . Specifically, given an initial control of the type $\frac{d\rho_0}{d\pi} \in L^{q'}(\pi)$ for some $q' < \infty$, it suffices to obtain a WPI with $\Phi(\cdot) = \|\cdot\|_{L^u(\pi)}^2$ (for mean-zero functions), where $u = \frac{2q'}{q}$. The following proposition, which is a consequence of a general L_p decay result due to [CCG12] and might be of independent interest, is crucial in the proof of Theorem 2. Specifically, it converts a WPI with $\Phi(\cdot) = \operatorname{Osc}(\cdot)^2$ into a different WPI with $\Phi(\cdot) = \|\cdot\|_{L^u(\pi)}^2$; thus, allows for a less restrictive initialization. We defer the proof of this proposition to Appendix A.2.

Proposition 3. Suppose π satisfies (WPI) with $\Phi(\cdot) = \operatorname{Osc}(f)^2$ and some $\beta_{WPI}(r)$. Then, for every u > 2, π also satisfies (WPI) with weighting β' and regularization function Φ' such that

$$\beta_{\mathsf{WPI}}'(r) = \beta_{\mathsf{WPI}}\left((r/5)^{\frac{u}{u-2}}\right) \ln\left((5/r)^{\frac{u}{u-2}}\right) \quad and \quad \Phi_u'(\cdot) = \left\|f - \mathbb{E}_{\pi}\left[f\right]\right\|_{L^u(\pi)}^2.$$
(2.4)

Additionally, provided that π does not satisfy a PI, π cannot satisfy a WPI with $\Phi = \Phi'_u$ for u = 2.

The last statement of the above proposition clarifies why we need to choose q' > q. In order to establish guarantees with q' = q, one would be required to obtain a WPI with $\Phi(\cdot) = \|\cdot\|_{L^2(\pi)}^2$, which is equivalent to a PI.

3 Langevin Monte Carlo for Heavy-Tailed Targets

In this section, we present our main convergence guarantee on LMC when the target satisfies (WPI) and ∇V is s-Hölder continuous for some $s \in (0, 1]$, that is,

$$\|\nabla V(x) - \nabla V(y)\| \le L \|x - y\|^s \quad \forall x, y \in \mathbb{R}^d.$$
 (s-Hölder)

The case s = 1 corresponds to the ubiquitous Lipschitz continuity where the target is smooth, and the regime where s < 1 is often termed as weak smoothness. Since our main focus is potentials that do not satisfy (PI), any order Hölder continuity is feasible.

Below, we state our main convergence result for a generic (WPI) with an unspecified β_{WPI} . We will explicitly derive its implications for specific targets in the subsequent sections.

Theorem 4. Suppose $\pi \propto e^{-V}$ satisfies (WPI) for some β_{WPI} and $\Phi(\cdot) = \text{Osc}(\cdot)^2$, and ∇V is (s-Hölder) continuous, with $\nabla V(0) = 0$ for simplicity. For any $q \in (1, \infty), q' \in (2q - 1, \infty]$ such that $R_{q'}(\mu_0 || \pi) < \infty$, define $\delta_0 \coloneqq \exp(qR_{q'}(\mu_0 || \pi)), m \coloneqq \frac{1}{2}\inf\{R : \pi(||x|| \ge R) \le \frac{1}{2}\}$, and

$$T := (2q-1) \left\{ \beta \left(\frac{1}{4\delta_0} \right) R_{2q-1}(\mu_0 \parallel \pi) + \beta \left(\frac{\varepsilon}{8\delta_0} \right) \ln \left(\frac{2}{\varepsilon} \right) \right\},$$

for $\varepsilon \leq q^{-1}$, where β is as in (2.2). Let $\hat{\pi}$ denote a modified version of π (explicitly defined in (A.3)) and assume, for simplicity, that ε^{-1} , $m, L, T, R_2(\mu_0 \parallel \hat{\pi}) \geq 1$. Then, for a sufficiently small step size h, denoting by μ_N , the law of the N-th iterate of LMC initialized at ρ_0 , after

$$N = \Theta_s \left(\frac{T^{1+1/s} dq^{1/s} L^{2/s}}{\varepsilon^{1/s}} \max\left\{ 1, \frac{\varepsilon^{1/(2s)} m^s}{L^{1/s-1} T^{1/(2s)} d}, \frac{\varepsilon^{1/(2s)} R_2(\mu_0 \parallel \hat{\pi})^{s/2}}{L^{1/s-1} T^{(1-s^2)/(2s)} d} \ln \left(\frac{qTLR_2(\mu_0 \parallel \hat{\pi})}{\varepsilon} \right)^{s/2} \right\} \right)$$

iterations of LMC, we obtain $R_q(\mu_N \parallel \pi) \leq \varepsilon$.

We make a few remarks. First, it is possible to find an initialization μ_0 , e.g. isotropic Gaussian, such that $R_q(\mu_0 \parallel \pi), R_q(\mu_0 \parallel \pi) \leq \tilde{\mathcal{O}}(d)$, with details provided in Lemma 30. In such a scenario, up to log factors, the last term in the maximum will never dominate. Additionally, the middle term will never dominate for sufficiently small ε , and in fact, in our model examples, it will not dominate even for $\varepsilon = 1$ due to proper control on m. Then our convergence rate reads

$$N = \tilde{\Theta}_s \left(\frac{q^{1+2/s} L^{2/s} d\left\{ d\beta(\frac{1}{4\delta_0}) + \beta(\frac{\varepsilon}{4\delta_0}) \right\}^{1+1/s}}{\varepsilon^{1/s}} \right)$$

In the case where π satisfies (PI), we set $\beta(r) = \beta_{PI}$ for any r > 0, and the above rate reduces to the rate implied by [CEL+22, Theorem 7].

The proof of Theorem 4 is given in Appendix A.3, and it is based on Theorem 2 and the Girsanov argument used in [CEL⁺22]. Two key distinctions are (i) our analysis is tailored for heavy-tailed targets, hence does not require any moment order to be defined, and also requires a finer control of different error terms to mitigate suboptimal rates, and (ii) it exploits the continuous time convergence under WPI, rather than better behaved stronger functional inequalities such as PI.

3.1 Examples

In this section, we focus on various heavy-tailed targets that do not satisfy (PI). In particular, we consider sampling from Cauchy-type measures in Section 3.1.2, which may not even have a defined expectation or any order moment for that matter; yet, we are able to provide convergence guarantees for LMC in Rényi divergence of any finite order.

3.1.1 Potentials with sub-linear tails

Consider the measure $\pi_{\alpha} \propto \exp(-V_{\alpha})$ where $V_{\alpha}(x) = (1 + ||x||^2)^{\alpha/2}$ with $\alpha \in (0, 1)$. This potential satisfies (*s*-Hölder) with s = 1 and L = 1. We analyze this potential as a substitute for $||x||^{\alpha}$ since the latter does not have continuous gradients, while the former still behaves similar to $||x||^{\alpha}$ for large ||x||. In the following Proposition, we present the β_{WPI} estimate for this potential.

Proposition 5. The measure $\pi_{\alpha}(x) \propto \exp(-V_{\alpha})$ with $\alpha \in (0,1)$ satisfies (WPI) with

$$\beta_{\mathsf{WPI}}(r) = \inf_{\gamma \in (0,2\alpha]} C_{\alpha} \left(d^{\frac{2(2-2\alpha+\gamma)}{\gamma}} + \ln(r^{-1})^{\frac{2-2\alpha+\gamma}{\alpha}} \right) \quad and \quad \Phi(\cdot) = \mathrm{Osc}(\cdot)^{2}, \tag{3.1}$$

where C_{α} is a constant depending only on α .

It is worth noting that this estimate is polynomial in dimension, improving the implicit and potentially exponential dependency implied by [CGGR10, Proposition 5.6] at the expense of introducing an additional

factor of $\ln(1/r)^2$ into the estimate. Specifically, they show that any potential of the form $V(x) = \psi(x)^{\alpha}$ for some convex non-negative ψ and $\alpha \in (0, 1)$ satisfies a WPI with $\beta_{\mathsf{WPI}}(r) = C_{\mathsf{WPI}}(1 + \ln(r^{-1})^{2(1/\alpha-1)})$ where $C_{\mathsf{WPI}} = C_{\mathsf{WPI}}(d, \alpha)$ is implicit in dimension, and in general it is not known how to achieve a dependency better than exponential via such techniques. Invoking Theorems 2 and 4 with the estimate for β_{WPI} given by Proposition 5, we can establish the following convergence guarantees for LMC and the Langevin diffusion.

Corollary 6. Consider the setting of Theorem 4 with the target measure $\pi_{\alpha} \propto \exp(-V_{\alpha})$. Denoting the distribution of the Langevin diffusion at time T with ρ_T , we have $R_q(\rho_T || \pi_{\alpha}) \leq \varepsilon$ whenever

$$T \ge C_{\alpha,q} \inf_{\gamma \in (0,2\alpha]} \left\{ \left(d^{\frac{2}{\gamma}(2-2\alpha+\gamma)} + R_{\infty}(\rho_0 \parallel \pi)^{\frac{2-2\alpha+\gamma}{\alpha}} \right) (R_q(\rho_0 \parallel \pi) + \ln(1/\varepsilon)) + \ln(1/\varepsilon)^{\frac{2-\alpha+\gamma}{\alpha}} \right\}$$

Further, denoting the distribution of LMC after N iterations with μ_N , we have $R_q(\mu_N \parallel \pi_\alpha) \leq \varepsilon$ if

$$N = \tilde{\Theta}_{\alpha,q} \left(\frac{(d^{4/\alpha} + R_{\infty}(\mu_0 \parallel \pi_{\alpha})^{4/\alpha}) R_q(\mu_0 \parallel \pi_{\alpha})^2 d}{\varepsilon} \max\left\{ 1, \frac{\sqrt{\varepsilon R_2(\mu_0 \parallel \hat{\pi}_{\alpha})}}{d} \right\} \right)$$

In particular, with $\rho_0 = \mu_0 = \mathcal{N}(0, I_d)$, and $\gamma = 2\alpha$, and using the bound on the initial Rényi divergence provided by Corollary 23 and Lemma 31, we obtain

$$T \gtrsim C_{\alpha,q} \left(d^{2/\alpha+1} + d^{2/\alpha} \ln(1/\varepsilon) + \ln(1/\varepsilon)^{2/\alpha+1} \right), \quad \text{and} \quad N = \tilde{\Theta}_{\alpha,q} \left(\frac{d^{4/\alpha+3}}{\varepsilon} \right),$$

where \gtrsim hides polylog factors in d.

As a final remark, if instead of the β_{WPI} estimate of Proposition 5, we use the estimate from [CGGR10] together with Theorem 4, we can obtain a rate for LMC which reads $N = \tilde{\Theta}_{\alpha,q} \left(C_{\mathsf{WPI}}^2 \frac{d^{4/\alpha-1}}{\varepsilon} \right)$. In particular, this can be seen as a smooth extrapolation of the rate given by [CEL+22] for the case of $\alpha \in [1,2]$. By showing that an LOI holds with some constant C_{LOI} , they obtain rates identical to ours once C_{WPI} is replaced by C_{LOI} , and the two rates match at $\alpha = 1$ as C_{LOI} and C_{WPI} will be equivalent up to an absolute constant in this case. However, when $\alpha < 1$, C_{WPI} of [CGGR10] has an implicit and potentially exponential dependence on dimension. Although our estimate in Proposition 5 improves this to polynomial, due to the additional log factor introduced into the bound, the rate will no longer smoothly extrapolate the rate of [CEL+22] to the regime $\alpha < 1$.

3.1.2 Potentials with logarithmic tails: Cauchy-type measures

In this section, we consider Cauchy-type measures of the form $\pi_{\nu} \propto \exp(-V_{\nu})$ where $V_{\nu}(x) = \frac{d+\nu}{2} \ln(1+||x||^2)$ with $\nu > 0$, which is Hölder continuous with s = 1 and $L \leq d + \nu$. In particular, π_{ν} belongs to the family of *d*-dimensional Student-t distributions with ν degrees of freedom; see e.g. [JR07, Kam18, HBE22, YLR22] for sampling from such distributions. Cauchy-type measures only have finite moments of order less than ν .

The following result presents a sharp estimate of β_{WPI} for π_{ν} .

Proposition 7. The measure $\pi_{\nu} \propto \exp(-V_{\nu})$ for $\nu > 0$ satisfies (WPI) with

$$\beta_{\mathsf{WPI}}(r) = \frac{2}{\nu} + 2\left(\frac{d}{\nu} + 1\right)r^{-2/\nu} \quad and \quad \Phi(\cdot) = \mathrm{Osc}(\cdot)^2.$$
(3.2)

Similar to the previous case, the above estimate improves the potentially exponential dimension dependence in [CGGR10, Proposition 5.4] to linear, while keeping r dependency the same. Employing the above estimate in Theorems 2 and 4, we obtain the following rate for LMC and the Langevin diffusion.

Corollary 8. Consider the setting of Theorem 4 with the target measure $\pi_{\nu} \propto \exp(-V_{\nu})$. Denoting the distribution of the Langevin diffusion at time T with ρ_T , we have $R_q(\rho_T || \pi_{\nu}) \leq \varepsilon$ whenever

$$T \ge C_{\nu}qd \exp\left(\frac{2qR_{\infty}(\rho_0 \parallel \pi)}{\nu}\right) \left(R_q(\rho_0 \parallel \pi) + (1/\varepsilon)^{2/\nu} \ln(1/\varepsilon)\right)$$

Furthermore, let $\nu \leq cd$ for some absolute constant c, and denote the distribution of LMC after N iterations with μ_N . Then, we have $R_q(\mu_N || \pi_{\nu}) \leq \varepsilon$ whenever

$$N = \tilde{\Theta}_{\nu,q} \left(\frac{d^5 e^{\frac{4(2q-1)}{\nu} R_{\infty}(\mu_0 \| \pi)}}{\varepsilon} \left(R_{2q-1}^2(\mu_0 \| \pi) + \left(\frac{1}{\varepsilon}\right)^{4/\nu} \right) \max\left\{ 1, \frac{\sqrt{\varepsilon R_2(\mu_0 \| \hat{\pi}) R_{\infty}(\mu_0 \| \pi)}}{d} \right\} \right)$$

We highlight that the dependence on the initial error is exponential in contrast to the polynomial dependence in the previous case. In other words, with a naive initialization, in the worst case, this dependency will negatively impact the complexity of sampling from Cauchy-type measures. However, if one initializes with an isotropic Gaussian with an appropriately scaled variance, using the estimates on the initial Rényi divergence in Corollary 22 and Lemma 31, we can obtain

$$T \ge C_{\nu,q} d^{2q/\nu+1} \left(\ln d + (1/\varepsilon)^{2/\nu} \ln(1/\varepsilon) \right) \quad \text{and} \quad N = \tilde{\Theta}_{\nu,q} \left(d^{\frac{4(2q-1)}{\nu}+5} / \varepsilon^{4/\nu+1} \right).$$

Finally, we remark that such an initialization may not be available in general; thus, the exponential dependence might be unavoidable.

4 Lower Bounds for LMC via Variance Decay

As demonstrated in the examples, the convergence guarantees given in Section 3 have worse dependence on the initial divergence when the tails of the target are heavier, with a sharp transition at (Cauchy-type) logarithmic tails, in which case the dependence on the initial error becomes exponential. In this section, we show that this is not due to a limitation in the analysis, but is in fact a property of LMC in heavy-tailed settings. We present a method for developing lower bounds for the convergence rate of LMC, with the goal of observing a similar dependence on the initial divergence and in particular, establishing that the exponential dependency for Cauchy-type measures is unavoidable in the worst case. First, we introduce the notation of complexity that we lower bound.

Definition 9. Let \mathcal{D} and \mathcal{D}' denote two divergences.

- The iteration complexity of LMC, denoted by $\mathscr{C}_{\mathcal{D},\mathcal{D}'}^{LMC}(\pi,h,\Delta_0,\varepsilon)$, is the smallest $k \in \mathbb{N}$ such for any μ_0 satisfying $\mathcal{D}'(\mu_0 \parallel \pi) \leq \Delta_0$, LMC with step-size h satisfies $\mathcal{D}(\mu_k \parallel \pi) \leq \varepsilon$.
- Similarly, the time complexity of the Langevin diffusion, $\mathscr{C}_{\mathcal{D},\mathcal{D}'}^{LD}(\pi, \Delta_0, \varepsilon)$, is the smallest $T \in \mathbb{R}_+$ such that for any ρ_0 satisfying $\mathcal{D}'(\rho_0 \| \pi) \leq \Delta_0$, the diffusion satisfies $\mathcal{D}(\rho_T \| \pi) \leq \varepsilon$.

For example, \mathcal{D} and \mathcal{D}' could be the KL divergence, the order-q Rényi divergence or the p-Wasserstein distance. By having this notion of complexity uniformly over all initial distributions whose initial error does not exceed Δ_0 , we mirror the types of bounds given in Section 3 and in the literature on LMC more broadly (see e.g. [CGLL22, Definition 5]). Indeed, in Section 3.1, we provided upper bounds for this quantity with $\mathcal{D} = R_q$ and $\mathcal{D}' = R_{q'}$.

Our methodology for developing lower bounds is based on the observation that in order for LMC to be close to the target in some divergence, its moments should match those of the target. We study the second moments of LMC and the Langevin diffusion and find that, by tracking their evolution through differential inequalities, we can relate their convergence to the convergence of processes that are more tractable. To relate this to the Rényi divergence, we rely on the strategy of variational representations; the most famous of these, the Donsker-Varadhan theorem, has the KL divergence represented as a supremum over a functional. We use a similar representation for the q-Rényi divergence with q > 1, which is due to [BDK⁺21], with a particular choice of test function to obtain

$$R_q(\rho \| \pi) \ge \ln \left(\rho(\| \cdot \|^2)^{\frac{q}{q-1}} / \pi(\| \cdot \|^{\frac{2q}{q-1}}) \right) \quad \text{whenever} \quad \pi(\| \cdot \|^{2q/(q-1)}) < \infty$$

Therefore, a strategy for obtaining lower bounds in the case where $\mathbb{E}||x_0||^2$ is large follows: (i) relate the decay of the Rényi divergence to the decay of $\mathbb{E}||x_k||^2$, (ii) lower bound the number of iterations required for $\mathbb{E}||x_k||^2$ to decay by the number of iterations for a more tractable process (e.g. gradient descent) to converge.

To strengthen the lower bounds obtained, we use a similar approach to develop an upper bound on the step size necessary to obtain a given accuracy. Details for the methodology can be found in Appendix B. We also remark that using variational representations of divergences to obtain lower bounds for statistical estimators goes back to the work of [CR51].

4.1 Heavy-Tailed Potentials and Slow Starts

In order to obtain lower bounds, we will make an assumption of the type

$$\|\nabla V(x)\| \lesssim \|x\|^{\alpha-1},$$

with $\alpha \in [0, 2]$, for sufficiently large ||x||, which will capture potentials that grow no faster than the order of $||x||^{\alpha}$ when $\alpha > 0$, and $\ln(||x||)$ when $\alpha = 0$. Under such a condition, we show that as $\alpha \to 0$, the dependence on initial error deteriorates. In particular, the dependence is logarithmic with $\alpha = 2$, becomes polynomial with $\alpha \in (0, 2)$, and finally turns exponential with $\alpha = 0$, as outlined below.

Theorem 10 (Three-step Phase Transition). Let $q \in (1, \infty)$, $q' \in [1, \infty]$ and the moment condition $\pi(\|\cdot\|^{2q/(q-1)}) < \infty$ holds. Suppose that

$$\|\nabla V(x)\| \le \frac{b\|x\|}{(1+\|x\|^2)^{1-\alpha/2}},\tag{4.1}$$

with $\alpha \in [0,2]$ and b > 0, and let $\nu := b - d$. Then for all sufficiently large Δ_0 (see (B.3), (B.4), and (B.5) for respective cases), the time complexity of the Langevin diffusion, and the iteration complexity of LMC, for obtaining an accuracy of 1 in q-Rényi divergence, is lower bounded as follows

$$\begin{split} \alpha &= 2: \qquad \mathscr{C}_{R_q,\mathsf{KL}}^{LMC}(\pi,h,\Delta_0,1) \gtrsim \frac{\ln(\Delta_0)}{h}, \qquad \mathscr{C}_{R_q,\mathsf{KL}}^{LD}(\pi,\Delta_0,1) \gtrsim \ln(\Delta_0), \\ \alpha &\in (0,2): \quad \mathscr{C}_{R_q,R_{q'}}^{LMC}(\pi,h,\Delta_0,1) \gtrsim \frac{d^{1-\alpha/2}\Delta_0^{\frac{(2-\alpha)^2}{2\alpha}}}{h}, \quad \mathscr{C}_{R_q,R_{q'}}^{LD}(\pi,\Delta_0,1) \gtrsim d^{1-\alpha/2}\Delta_0^{\frac{(2-\alpha)^2}{2\alpha}}, \\ \alpha &= 0: \qquad \mathscr{C}_{R_q,R_{q'}}^{LMC}(\pi,h,\Delta_0,1) \gtrsim \frac{d}{\nu h} \exp(\Delta_0/\nu), \quad \mathscr{C}_{R_q,R_{q'}}^{LD}(\pi,\Delta_0,1) \gtrsim \frac{d}{\nu} \exp(\Delta_0/\nu), \end{split}$$

where \gtrsim hides a constant depending only on b and α . The lower bounds hold for all h > 0, except when $\alpha = 2$ where they hold for $h < b^{-1}$.

In Theorem 10, the case of $\alpha = 2$ corresponds to tails no lighter than Gaussian and indeed, the lower bounds that we obtain reproduce the dependence of Δ_0 known in the Gaussian setting ([VW19]). The case of $\alpha \in (0, 2)$ corresponds to potentials with tail growth similar to $||x||^{\alpha}$ and the case of $\alpha = 0$ corresponds to Cauchy-type logarithmic tails. Indeed, the generalized Cauchy potentials V_{ν} and the sub-linear potentials V_{α} from Section 3.1 satisfy (4.1) with $b = d + \nu$ and $b = \alpha$ respectively. For these examples, our lower bounds recover the polynomial and exponential dependence on the initial error given in the upper bounds of Section 3. More generally, we point out that (4.1) is satisfied for any smooth ∇V with $\nabla V(0) = 0$ and $||\nabla V(x)|| \leq ||x||^{\alpha-1}$ for large ||x||. Similar to the upper bounds, the implicit assumption $\nabla V(0) = 0$ is made only for the simplicity of presentation. In fact, to achieve the most generality, even this assumption can be relaxed to a bound of the type $\langle \nabla V(x), x \rangle \leq ||x||^{\alpha}$ for large ||x|| while recovering similar results. One can interpret such a condition as *reversed dissipativity*. While dissipativity is used in the literature to ensure a lower bound on tail growth for obtaining upper bounds on moments (see e.g. [RRT17, EMS18, EHZ22, FR21, YBVE21]), the reversed condition imposes an upper bound on the tail growth that leads to obtaining lower bounds on the moments.

The above result highlights one reason why LMC can be seen to perform worse in heavy-tailed settings. As the tail becomes heavier, LMC exhibits a slow start behavior by having a worse dependence on the initial divergence. Showing that this is a property of the Langevin diffusion and not due to discretization motivates using alternative diffusions for sampling from heavy-tailed targets (see, for example, [LWME19, §ZTG20, HBE22, ZZ22, CDSX22]).

To complete our exposition, we note that while the results of Theorem 10 are stated for any fixed choice of step size, in practice we have to ensure step size is small enough for the discretization to not harm the convergence of LMC. This usually leads to additional dependence on dimension or final accuracy. In a special case where the target potential is radially symmetric, the following proposition provides a general tool for determining suitable ranges of step size for LMC, hence completing the complexity lower bound in conjunction with Theorem 10.

Proposition 11 (Step-size upper bound). Suppose the potential is radially symmetric with $V(x) = f(||x||^2)$ and $g: \mathbb{R}_+ \to \mathbb{R}_+$ given by $g(r) = (1 - 2hf'(r))^2 r$ is convex and non-decreasing. Let $\varepsilon > 0, q > 1$. Suppose that $\mathbb{E}[||x_0||^2] > \sigma_{\varepsilon}^2$ where $\sigma_{\varepsilon}^2 = e^{\frac{q-1}{q}\varepsilon}\pi(||\cdot||^{\frac{2q}{q-1}})^{\frac{q-1}{q}}$. Then it must hold that,

$$\inf_{k\in\mathbb{N}} R_q(\rho_k \| \pi) < \varepsilon \implies h \le \frac{1}{f'(\sigma_{\varepsilon}^2)} \left(1 - \frac{d}{2f'(\sigma_{\varepsilon}^2)\sigma_{\varepsilon}^2} \right).$$

We conclude our discussion by considering an example that demonstrates an application of the main tools developed in this paper. Specifically, we consider the setting of Section 3.1.2 and recall the Cauchy-type target, π_{ν} . The following corollary, which gives a sharp characterization of the convergence behavior, is a direct consequence of the lower bound from Theorem 10 and the upper bound from Corollary 8.

Corollary 12. Let q > 1 and suppose that $\nu > \frac{2q}{q-1}$. Then for any $\Delta_0 \ge C_q(1 + \nu \ln(d + \nu))$, we have the following bound for the Langevin diffusion,

$$\frac{d}{4\nu} \exp\left(\frac{\Delta_0}{\nu}\right) \le \mathscr{C}_{R_q,R_\infty}^{LD}(\pi,\Delta_0,1) \le \frac{d}{\nu} \exp\left(\frac{C_q \Delta_0}{\nu}\right)$$

Further, LMC with an appropriate step size h > 0 satisfies

$$\frac{d}{4h\nu} \exp\left(\frac{\Delta_0}{\nu}\right) \le \mathscr{C}_{R_q,R_\infty}^{LMC}(\pi,h,\Delta_0,1) \le \frac{d}{h\nu} \exp\left(\frac{C_q\Delta_0}{\nu}\right)$$

The above result states that the exponential dependence on the initial error for LMC and the diffusion is unavoidable, unless there is a good initialization available. Notice that, by Proposition 11, to obtain $R_q(\mu_N \parallel \pi) \leq \varepsilon$, step size needs to be sufficiently small. Incorporating this into the above bound, we obtain the following lower bound for the iteration complexity of LMC when $d \geq \nu - p$:

$$\mathscr{C}_{R_q,R_{q'}}^{\mathrm{LMC}}(\pi,h,\Delta_0,\varepsilon) \geq \frac{(d+\nu)^2}{48e\nu} \min\left\{\varepsilon^{-1},\frac{\nu-p}{p},\frac{(\nu-p)d}{(2+p)\nu}\right\}\exp\left(\frac{\Delta_0}{\nu}\right),$$

where $p = \frac{2q}{q-1}$. Note that the step-size bound leads to additional dependence on d and ε . However, the dependence on ε only reveals itself when ν and d are relatively large.

5 Conclusion

We provided convergence guarantees for LMC and Langevin diffusion, for target distributions $\pi \propto e^{-V}$ satisfying a WPI. This corresponds to potentials with tails that behave like $V(x) \sim ||x||^{\alpha}$ for $\alpha \in (0, 1]$, and also covers Cauchy-type densities in the case $\alpha \downarrow 0$. Our results push the program initiated by [VW19] to its limits; by providing explicit WPI constants for specific examples, we obtained guarantees demonstrating that targets with heavier tails lead to a worse dependence on the initial error. Particularly, the dependence on initial error is a polynomial of order $\frac{(2-\alpha)^2}{2\alpha}$ for $\alpha > 0$, with a phase transition at $\alpha = 0$, i.e. Cauchy-type logarithmic tails, for which the dependence becomes exponential. Additionally, we established lower bounds under generic tail growth conditions that asserted such dependence on the initial error is unavoidable unless suitable warm start initializations are available.

Our quantitative results (via upper and lower bounds) highlight the precise limitations of LMC for heavytailed sampling and provide further motivations to develop a complete complexity theory of heavy-tailed sampling by discretizing other diffusions like specific Itô diffusions or stable-driven diffusions, an area of research which is still in its infancy.

One limitation of our upper bounds is the fact that the step size needs to be chosen in advance according to the number of iterations, and with a fixed step size, the upper bounds diverge as $N \to \infty$. However, as pointed out by [CEL⁺22], this is a general limitation of any analysis that does not assume the log-Sobolev inequality. We leave the stability of fixed step size LMC in the number of iterations under heavy-tailed targets as an open direction for future research.

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A Proofs of Section 3

We begin by presenting the proof of Theorem 2 by obtaining WPIs with appropriate Φ' corresponding to $R_{q'}(\rho_0 \parallel \pi)$ using the canonical WPI with $\Phi = \text{Osc}(\cdot)^2$, and establishing a differential inequality for Rényi divergence along the Langevin diffusion using the aforementioned WPIs.

A.1 Proof of Theorem 2

Following [VW19], define the following quantities

$$F_q(\rho \parallel \pi) \coloneqq \mathbb{E}_{\pi} \left[\left(\frac{\rho}{\pi} \right)^q \right], \qquad G_q(\rho \parallel \pi) \coloneqq \mathbb{E}_{\pi} \left[\left(\frac{\rho}{\pi} \right)^q \left\| \nabla \log \frac{\rho}{\pi} \right\|^2 \right] = \frac{4}{q^2} \mathbb{E}_{\pi} \left[\left\| \nabla \left(\frac{\rho}{\pi} \right)^{q/2} \right\|^2 \right].$$

Then [VW19, Lemma 6] shows that along the Langevin diffusion

$$\frac{\partial R_q(\rho_t \parallel \pi)}{\partial t} = -\frac{qG_q(\rho_t \parallel \pi)}{F_q(\rho_t \parallel \pi)}.$$
(A.1)

Furthermore, we have the following estimates on the quantities appearing in our functional inequalities.

Lemma 13 ([VW19]). Let $f = \left(\frac{\rho_t}{\pi}\right)^{q/2}$. Then,

$$\mathbb{E}_{\pi}\left[\left\|\nabla f\right\|^{2}\right] = \frac{q^{2}}{4}G_{q}(\rho_{t} \parallel \pi),$$

$$\operatorname{Var}_{\pi}(f) \geq F_{q}(\rho_{t} \parallel \pi)(1 - e^{-R_{q}(\rho_{t} \parallel \pi)}).$$

Proof. The equality follows by definition. For the inequality,

$$\begin{aligned} \operatorname{Var}_{\pi}(f) &= F_{q}(\rho_{t} \parallel \pi) - F_{q/2}(\rho_{t} \parallel \pi)^{2} \\ &= e^{(q-1)R_{q}(\rho_{t} \parallel \pi)} - e^{(q-2)R_{q/2}(\rho_{t} \parallel \pi)} \\ &\geq e^{(q-1)R_{q}(\rho_{t} \parallel \pi)} - e^{(q-2)R_{q}(\rho_{t} \parallel \pi)} \\ &= F_{q}(\rho_{t} \parallel \pi) \Big(1 - e^{-R_{q}(\rho_{t} \parallel \pi)} \Big), \end{aligned}$$

where we used the fact that $q \mapsto R_q(\rho_t || \pi)$ is non-decreasing.

Proof. [Proof of Theorem 2.] Suppose π satisfies a Poincaré inequality with β_{WPI} and Φ . Assume $\Phi\left(\left(\frac{\rho_t}{\pi}\right)^{q/2}\right) \leq \delta_0$ for all $t \in \mathbb{R}_+$. Choose $f = \left(\frac{\rho_t}{\pi}\right)^{q/2}$ in (WPI). Lemma 13 yields

$$\begin{aligned} \frac{q^2 G_q(\rho_t \parallel \pi)}{4 F_q(\rho_t \parallel \pi)} &\geq \frac{1 - e^{-R_q(\rho_t \parallel \pi)}}{\beta_{\mathsf{WPI}}(r)} - \frac{r \Phi(f)}{\beta_{\mathsf{WPI}}(r) F_q(\rho_t \parallel \pi)} \\ &\geq \frac{1 - e^{-R_q(\rho_t \parallel \pi)}}{\beta_{\mathsf{WPI}}(r)} - \frac{r \delta_0}{\beta_{\mathsf{WPI}}(r)}. \end{aligned}$$

Hence

$$\frac{\partial R_q(\rho_t \parallel \pi)}{\partial t} \leq \frac{-4(1 - e^{-R_q(\rho_t \parallel \pi)}) + 4r\delta_0}{q\beta_{\mathsf{WPI}}(r)}$$

Thus with $R_q(\rho_t \parallel \pi) \ge 1$ we have

$$\frac{\partial R_q(\rho_t \parallel \pi)}{\partial t} \le \frac{-2 + 4r\delta_0}{q\beta_{\mathsf{WPI}}(r)}$$

and for $R_q(\rho_t \parallel \pi) < 1$ we have

$$\frac{\partial R_q(\rho_t \parallel \pi)}{\partial t} \le \frac{-2R_q(\rho_t \parallel \pi) + 4r\delta_0}{q\beta_{\mathsf{WPI}}(r)},$$

Integration and Grönwall's lemma yield

$$R_q(\rho_t \parallel \pi) \leq \begin{cases} R_q(\rho_0 \parallel \pi) - \frac{2 - 4r\delta_0}{\beta_{\mathsf{WPI}}(r)q}t & \text{if } R_q(\rho_0 \parallel \pi), R_q(\rho_t \parallel \pi) \geq 1\\ e^{-\frac{2t}{\beta_{\mathsf{WPI}}(r)q}} (R_q(\rho_0 \parallel \pi) - 2r\delta_0) + 2r\delta_0 & \text{if } R_q(\rho_0 \parallel \pi) < 1 \end{cases}$$

Suppose $q' = \infty$. In this case, we can choose the original WPI with $\Phi(\cdot) = Osc(\cdot)^2$. Then, we need to show $\operatorname{Osc}\left(\left(\frac{\rho_t}{\pi}\right)^{q/2}\right)^2 \leq \left\|\frac{\rho_0}{\pi}\right\|_{L^{\infty}(\pi)}^q$. Notice that $P_t \frac{\rho_0}{\pi} = \frac{\rho_t}{\pi}$ (which one can verify by the tower property of the conditional expectation and the symmetry of the Markov process). Then we have,

$$\operatorname{Osc}\left(\left(P_t \frac{\rho_0}{\pi}\right)^{q/2}\right)^2 \le \left\|P_t \frac{\rho_0}{\pi}\right\|_{L^{\infty}(\pi)}^q \le \left\|\frac{\rho_0}{\pi}\right\|_{L^{\infty}(\pi)}^q,$$

where the second inequality follows from the contraction property of the semigroup. Thus, we can set $\delta_0 = \left\| \frac{\rho_0}{\pi} \right\|_{L^{\infty}(\pi)}^q = \exp(qR_{\infty}(\rho_0 \parallel \pi)),$ which completes the proof of the $q' = \infty$ case.

Now, suppose $q < q' < \infty$. Using $u = \frac{2q'}{q}$ in Proportion 3, π satisfies a WPI with $\beta'_{WPI} = \beta$ and $\Phi'(\cdot) = \|\cdot - \mathbb{E}_{\pi}[\cdot]\|^2_{L^{2q'/q}(\pi)}$, where β is given in the statement of Theorem 2. Additionally, by Lemma 33, we have

$$\Phi'\left(\left(\frac{\rho_t}{\pi}\right)^{q/2}\right) \le \left\|\left(\frac{\rho_t}{\pi}\right)^{q/2}\right\|_{L^{2q'/q}(\pi)}^2 = \left\|P_t\frac{\rho_0}{\pi}\right\|_{L^{q'}(\pi)}^q \le \left\|\frac{\rho_0}{\pi}\right\|_{L^{q'}(\pi)}^q \le \exp(qR_{q'}(\rho_0 \parallel \pi)).$$
is case, we can choose $\delta_0 = \exp(qR_{q'}(\rho_0 \parallel \pi))$, which finishes the proof.

Thus in this case, we can choose $\delta_0 = \exp(qR_{q'}(\rho_0 \parallel \pi))$, which finishes the proof.

Proof of Proposition 3 A.2

By Theorem 2.3 of [RW01], in order to show (2.4) it suffices to show that for every mean-zero smooth f

$$\|P_t f\|_{L^2(\pi)}^2 \le \xi(t) \|f\|_{L^u(\pi)}^2 \tag{A.2}$$

such that $\lim_{t\to\infty} \xi(t) = 0$. Then, (2.4) holds with

$$\beta'_{\mathsf{WPI}}(r) = 2r \inf_{s} \frac{1}{s} \xi^{-1}(s \exp(1 - s/r)).$$

Furthermore, by Lemma 5.1 of [CCG12], we have

$$\|P_t f\|_{L^2(\pi)}^2 \le 4^{1+2/u} \|f\|_{L^u(\pi)}^2 \left(\sup_{\{g: \mathbb{E}_\mu[g]=0\}} \frac{\|P_t g\|_{L^2(\pi)}^2}{\|g\|_{\infty}^2} \right)^{1-2/u}.$$

Moreover, by Theorem 2.1 of [RW01], for any smooth bounded g we have

$$||P_t g||^2_{L^2(\pi)} \le 5\lambda(t) ||g||^2_{\infty},$$

where

$$\lambda(t) = \inf\{r > 0 : \frac{1}{2}\beta_{\mathsf{WPI}}(r)\log(1/r) \le t\}.$$

Hence, (A.2) holds with $\xi(t) = 5\lambda(t)^{1-2/u}$. From this, we observe that

$$\beta_{\mathsf{WPI}}'(r) \le 2\xi^{-1}(r) = 2\lambda^{-1} \Big((r/5)^{\frac{u}{u-2}} \Big) = \beta_{\mathsf{WPI}} \Big((r/5)^{\frac{u}{u-2}} \Big) \log \Big((5/r)^{\frac{u}{u-2}} \Big).$$

Finally, suppose π does not satisfy a Poincaré inequality, but satisfies (2.4) for some β'_{WPI} with u = 2. Then, for any r > 0 we have

$$(1-r)\operatorname{Var}(f) \leq \beta_{\mathsf{WPI}}(r)\mathbb{E}\Big[\|\nabla f\|^2\Big].$$

Thus, π satisfies a Poincaré inequality (with a constant at most $2\beta_{WPI}(1/2)$), which is a contradiction.

A.3 Proof of Theorem 4

First, we present the following lemma which enables us to control the discretization error $R_{2q}(\mu_N \parallel \rho_T)$. This proposition can be retrieved by a careful evaluation of the terms in the proof of Proposition 25 of [CEL+22]. We avoid the simplifications made by [CEL+22] since that can affect our rate due to $m = \mathcal{O}(d^{1/\alpha})$ for π_{α} defined in Section 3 with $\alpha \in (0, 1)$, and also T can be exponential in d for specific examples.

Proposition 14 ([CEL⁺22, Proposition 25]). Suppose ∇V is s-Hölder continuous with constant L and $\nabla V(0) = 0$. Let $m \coloneqq \frac{1}{2} \inf\{R : \pi(||x|| \ge R) \le \frac{1}{2}\}$. Define

$$\hat{\pi}(x) \propto \exp\left(-V(x) - \frac{1}{6144T} \max\{\|x\| - 2m, 0\}^2\right).$$
 (A.3)

Assume for simplicity ε^{-1} , $m, L, T, R_2(\rho_0 \parallel \hat{\pi}) \ge 1$, Then, for $q \le \varepsilon^{-1}$, if the step size satisfies

$$h \le \mathcal{O}_s\left(\frac{\varepsilon^{1/s}}{dq^{1/s}L^{2/s}T^{1/s}}\min\left\{1, \frac{L^{1/s-1}T^{1/(2s)}d}{\varepsilon^{1/(2s)}m^s}, \frac{L^{1/s-1}T^{(1-s^2)/(2s)}d}{\varepsilon^{1/(2s)}R_2(\rho_0 \parallel \hat{\pi})^{s/2}\ln(N)^{s/2}}\right\}\right)$$

we have for T = Nh

$$R_q(\mu_N \parallel \rho_T) \le \varepsilon.$$

Remark. If the first moment of π is finite, we have $m \leq \mathbb{E}_{\pi}[||x||]$ by Markov's inequality. In fact, the original result in [CEL⁺22] is presented using $m = \mathbb{E}_{\pi}[||x||]$. However, the result is still valid with the choice of m in Proposition 14, which is useful for targets with infinite first moment. With this proposition in hand, we are ready to state the proof of Theorem 4.

Proof. [Proof of Theorem 4] By Theorem 2 we have $R_{2q-1}(\rho_T || \pi) \leq \frac{\varepsilon}{2}$. Furthermore, by choosing the step-size as in Proposition 14, we have $R_{2q}(\mu_N || \rho_T) \leq \frac{\varepsilon}{2}$. By the weak triangle inequality of Rényi divergence, we have $R_q(\mu_N || \pi) \leq \varepsilon$ for

$$N = \frac{T}{h} = \Theta_s \bigg(\frac{dq^{1/s} L^{2/s} T^{1+1/s}}{\varepsilon^{1/s}} \max\bigg\{ 1, \frac{\varepsilon^{1/(2s)} m^s}{L^{1/s-1} T^{1/(2s)} d}, \frac{\varepsilon^{1/(2s)} R_2(\rho_0 \parallel \hat{\pi})^{s/2} \ln(N)^{s/2}}{L^{1/s-1} T^{(1-s^2)/(2s)} d} \bigg\} \bigg).$$

A further simplification of this result yields the statement of the Theorem.

A.4 Computing Weak Poincaré Constants

In this section, we will provide WPI estimates for our model examples $\pi_{\nu}(x) \propto (1 + ||x||^2)^{\frac{d+\nu}{2}}$ and $\pi_{\alpha}(x) \propto \exp(-(1 + ||x||^2)^{\alpha/2})$.

We will use the following chain of implications to establish WPIs with suitable dimension dependencies:

Weighted Poincaré \Rightarrow Converse Poincaré \Rightarrow Weak Poincaré

In particular, to obtain a WPI from a converse Poinacré inequality, we can use the following result due to [CGGR10].

Lemma 15 (Theorem 5.1 of [CGGR10]). Assume π satisfies a converse Poincaré inequality, i.e.

$$\inf_{c} \int |f(x) - c|^2 w(x) d\pi(x) \le C \int \left\| \nabla f(x) \right\|^2 d\pi(x)$$

for some non-negative weight function w, such that $\int w d\pi < \infty$. Define $G(r) := \inf\{u : \pi(w \le u) > r\}$. Then, π satisfies a WPI with $\Phi(\cdot) = \operatorname{Osc}(\cdot)^2$ and $\beta_{\mathsf{WPI}}(r) = \frac{C}{G(r)}$.

Hence, our main effort is to establish a converse Poincaré inequality for our model examples. In fact, for generalized Cauchy measures, we can immediately use a result of [BL09].

Lemma 16 (Corollary 3.2 of [BL09]). Let $\pi_{\nu}(x) \propto (1 + ||x||^2)^{-(d+\nu)/2}$ with p > 0. Then, for any smooth f, we have the following converse (weighted) Poincaré inequality

$$\inf_{c} \int |f(x) - c|^{2} w(x) d\pi_{\nu}(x) \le C_{d,\nu} \int ||\nabla f(x)||^{2} d\pi_{\nu}(x), \tag{A.4}$$

for $w(x) = \frac{1}{1 + \|x\|^2}$, with

$$C_{d,\nu} \coloneqq \begin{cases} \frac{1}{d+\nu} & if\nu \ge d+2\\ \frac{2}{\nu} & otherwise \end{cases}$$

Thus, along with a concentration bound, we can invoke Lemma 15 to estimate β_{WPI} for generalized Cauchy measures.

Proof. [Proof of Proposition 7] In order to invoke Lemma 15, we need to estimate $\pi(w \le u)$ for the weight function given by Lemma 16 $w(x) = (1 + ||x||^2)^{-1}$. By Lemma 29, we have

$$\pi(w \le u) = \pi \left(\|x\| \ge \sqrt{u^{-1} - 1} \right) \le (d + \nu)^{\nu/2} (u^{-1} - 1)^{-\nu/2}.$$

Choosing $u^{-1} = 1 + (d + \nu)r^{-2/\nu}$, we obtain by invoking Lemma 15

$$\beta_{\text{WPI}}(r) = \frac{2}{\nu} + 2\left(\frac{d}{\nu} + 1\right)r^{-2/\nu}.$$

Estimating β_{WPI} for π_{α} is more involved as we do not readily have a suitable converse Poincaré inequality. We will work towards this by first deriving a weighted Poincaré inequality using the perturbation argument.

Lemma 17. Let $\pi_{\alpha}(x) \propto \exp\left(-\left(1+\|x\|^2\right)^{\alpha/2}\right)$ with $\alpha \in (0,1)$. Then for any smooth f, we have the following weighted Poincaré inequality

$$\operatorname{Var}_{\pi_{\alpha}}(f) \le eC_{d,\alpha} \int w(x)^2 \|\nabla f(x)\|^2 d\pi_{\alpha}(x).$$
(A.5)

for $w(x) = ||x||^{1-\alpha}$ with $C_{d,\alpha}$ satisfying (A.6).

Proof. Let $\tilde{\pi}_{\alpha}(x) \propto \exp(-\|x\|^{\alpha})$. According to [CGGR10, Proposition 4.7], $\tilde{\pi}_{\alpha}$ satisfies a weighted Poincaré inequality with weight $w(x)^2 = \|x\|^{2(1-\alpha)}$ and parameter $C_{d,\alpha}$ satisfies

$$\frac{d}{\alpha^3} \le C_{d,\alpha} \le 12\frac{d}{\alpha^3} + \frac{d+\alpha}{\alpha^4} \tag{A.6}$$

Meanwhile, $\frac{d\pi_{\alpha}}{d\tilde{\pi}_{\alpha}} = \exp(k(x))$ with $k(x) = -\left(1 + \|x\|^2\right)^{\alpha/2} + \|x\|^{\alpha} + \text{constant. Notice that Osc}(k) = 1$. Then, (A.5) follows from the perturbation property of the weighted Poincaré inequality.

In the next step, by an argument similar to [BL09, Proposition 3.3], we transform the weighted Poincaré inequality for π_{α} to a converse Poincaré inequality.

Lemma 18. Let $\pi_{\alpha}(x) \propto \exp\left(-\left(1+\|x\|^2\right)^{\alpha/2}\right)$ with $\alpha \in (0,1)$. Let $\gamma \in (0,2\alpha]$. Then for any smooth g, we have the following converse weighted Poincaré inequality

$$\inf_{c} \int \frac{|f(x) - c|^{2}}{\left(a + b \|x\|^{2}\right)^{1 - \alpha + \gamma/2}} d\pi_{\alpha}(x) \leq \left[1 - (1 - \alpha + \gamma/2)b^{\frac{1}{2}}a^{-\frac{1}{2}(\alpha - \gamma/2)}\right]^{-2} \int \|\nabla f(x)\|^{2} d\pi_{\alpha}(x)$$
(A.7)

with $a = \frac{\gamma(eC_{d,\alpha})^{\frac{2}{\gamma}}}{2(1-\alpha)+\gamma}$, $b = \frac{2(1-\alpha)}{2(1-\alpha)+\gamma}$, and $C_{d,\alpha}$ satisfies (A.6)

Proof. First we apply Young's inequality to bound the weights in Lemma 17.

$$eC_{d,\alpha} \|x\|^{2(1-\alpha)} = \left[(eC_{d,\alpha})^{\frac{2}{2(1-\alpha)+\gamma}} \|x\|^{\frac{4(1-\alpha)}{2(1-\alpha)+\gamma}} \right]^{1-\alpha+\gamma/2}$$
$$\leq \left[\frac{(eC_{d,\alpha})^{\frac{2}{\gamma}}}{\frac{2(1-\alpha)+\gamma}{\gamma}} + \frac{\|x\|^{2}}{\frac{2(1-\alpha)+\gamma}{2(1-\alpha)}} \right]^{1-\alpha+\gamma/2}$$
$$=: \left(a + b\|x\|^{2} \right)^{1-\alpha+\gamma/2}$$

Therefore, Lemma 17 yields

$$\operatorname{Var}_{\pi_{\alpha}}(g) \leq \int \|\nabla g(x)\|^2 (a+b\|x\|^2)^{1-\alpha+\gamma/2} d\pi_{\alpha}(x)$$
 (A.8)

The rest of the proof proceeds similarly to that of [BL09, Proposition 3.3]. Consider g(x) = (f(x) - f(x)) $c)\left(a+b\|x\|^2\right)^{-\frac{1-\alpha+\gamma/2}{2}}$, where c is chosen such that g has mean 0. Then we have

$$\nabla g(x) = \nabla f(x) \left(a + b \|x\|^2 \right)^{-\frac{1-\alpha+\gamma/2}{2}} - b(1-\alpha+\gamma/2) \left(a + b \|x\|^2 \right)^{-\frac{1-\alpha+\gamma/2}{2}-1} (f(x)-c)x$$

By the elementary inequality $||u+v||^2 \leq \frac{r}{r-1} ||u||^2 + r ||v||^2$ for any r > 1, we have

$$\begin{aligned} \|\nabla g(x)\|^{2} &\leq \frac{r}{r-1} \|\nabla f(x)\|^{2} \left(a+b\|x\|^{2}\right)^{-(1-\alpha+\gamma/2)} \\ &+ rb^{2}(1-\alpha+\gamma/2)^{2} \frac{\|x\|^{2}}{a+b\|x\|^{2}} \frac{1}{\left(a+b\|x\|^{2}\right)^{\alpha-\gamma/2}} \left(a+b\|x\|^{2}\right)^{-2(1-\alpha+\gamma/2)} (f(x)-c)^{2} \\ &\leq \frac{r}{r-1} \|\nabla f(x)\|^{2} \left(a+b\|x\|^{2}\right)^{-(1-\alpha+\gamma/2)} \\ &+ rb^{2}(1-\alpha+\gamma/2)^{2} b^{-1} a^{-\alpha+\gamma/2} \left(a+b\|x\|^{2}\right)^{-2(1-\alpha+\gamma/2)} (f(x)-c)^{2} \end{aligned}$$

Applying (A.8) to g, we obtain

$$\int (f(x) - c)^2 \left(a + b \|x\|^2\right)^{-(1 - \alpha + \gamma/2)} d\pi_\alpha(x) \le \frac{r}{r - 1} \int \|\nabla f(x)\|^2 d\pi_\alpha(x) + r b a^{-\alpha + \gamma/2} (1 - \alpha + \gamma/2)^2 \int (f(x) - c)^2 \left(a + b \|x\|^2\right)^{-(1 - \alpha + \gamma/2)} d\pi_\alpha(x)$$

Since $\gamma \in (0, 2\alpha]$, $b \in (0, 1)$, a > 1 and $1 - \alpha + \gamma/2 \in (1 - \alpha, 1)$, we have

$$\int f(x)^2 \left(a+b\|x\|^2\right)^{-(1-\alpha+\gamma/2)} d\pi_{\alpha}(x) \le \frac{r}{r-1} \frac{1}{1-rba^{-\alpha+\gamma/2}(1-\alpha+\gamma/2)^2} \int \|\nabla f(x)\|^2 d\pi_{\alpha}(x)$$
(A.7) follows by choosing the optimal $r^* = b^{-\frac{1}{2}} a^{\frac{1}{2}(\alpha-\gamma/2)}(1-\alpha+\gamma/2)^{-1} > 1.$

Last, (A.7) follows by choosing the optimal $r^* = b^{-\frac{1}{2}} a^{\frac{1}{2}(\alpha - \gamma/2)} (1 - \alpha + \gamma/2)^{-1} > 1$.

With this converse Poincaré inequality, we are ready to invoke Lemma 15 and obtain a WPI for π_{α} . **Proof.** [Proof of Proposition 5] Let a and b be defined according to the statement of Lemma 18. In order to invoke Lemma 15, we need to estimate $\pi_{\alpha}(w \le u)$ for the weight $w(x) = \left(a + b \|x\|^2\right)^{-(1-\alpha+\gamma/2)}$. Using the tail bound of Lemma 28. f tail bound of Lemma 28, f

$$\pi_{\alpha}(w(x) \le u) = \pi_{\alpha} \left(\|x\| \ge b^{-\frac{1}{2}} \sqrt{u^{-\frac{1}{1-\alpha+\gamma/2}} - a} \right)$$
$$\le e^{\frac{1}{2}} 2^{d/\alpha} \exp\left(-\frac{1}{2} \left(1 + b^{-1} \left(u^{-\frac{1}{1-\alpha+\gamma/2}} - a\right)\right)^{\alpha/2}\right).$$

Choosing u such that the above is at most r, we obtain

$$G(r) := \inf\{u : \pi_{\alpha}(w(x) \le u) > r\}$$

$$\geq \left\{a + b\left[\left(1 + \frac{2d}{\alpha}\ln 2 + 2\ln(r^{-1})\right)^{\frac{2}{\alpha}} - 1\right]\right\}^{-(1-\alpha+\gamma/2)}$$

$$\geq \left\{a + b\left(1 + \frac{2d}{\alpha}\ln 2 + 2\ln(r^{-1})\right)^{\frac{2}{\alpha}}\right\}^{-(1-\alpha+\gamma/2)}.$$

Therefore, by Lemma 15 the weak Poincaré constant satisfies

$$\begin{aligned} \alpha_{\mathsf{WPI}}(r) &\leq \left[1 - (1 - \alpha + \gamma/2)b^{\frac{1}{2}}a^{-\frac{1}{2}(\alpha - \gamma/2)}\right]^{-2} \left\{a + b\left(1 + \frac{2d}{\alpha}\ln 2 + 2\ln(r^{-1})\right)^{\frac{2}{\alpha}}\right\}^{1 - \alpha + \gamma/2} \\ &\leq \left(1 - a^{-\frac{1}{2}(\alpha - \gamma/2)}\right)^{-2} \left(a^{1 - \alpha + \gamma/2} + \left(1 + \frac{2d}{\alpha}\ln 2 + 2\ln(r^{-1})\right)^{\frac{2(1 - \alpha) + \gamma}{\alpha}}\right) \end{aligned}$$

Recall $a = \frac{\gamma}{2(1-\alpha)+\gamma} (eC_{d,\alpha})^{\frac{2}{\gamma}} < (eC_{d,\alpha})^{\frac{2}{\gamma}}, \gamma \in (0, 2\alpha]$, and notice that $\inf_{\gamma} a > 1$. Therefore,

$$\alpha_{\mathsf{WPI}}(r) \leq \frac{3^{\frac{2-3\alpha+\gamma}{\alpha}} \vee 1}{\left(1 - a^{-\frac{1}{2}(\alpha-\gamma/2)}\right)^2} \left(\left(eC_{d,\alpha}\right)^{\frac{2(2-2\alpha+\gamma)}{\gamma}} + 1 + \left(\frac{2\ln 2}{\alpha}\right)^{\frac{2-2\alpha+\gamma}{\alpha}} d^{\frac{2-2\alpha+\gamma}{\alpha}} + 2^{\frac{2-2\alpha+\gamma}{\alpha}} \ln\left(r^{-1}\right)^{\frac{2-2\alpha+\gamma}{\alpha}} \right)$$
$$\leq C_{\alpha} \left(d^{\frac{2(2-2\alpha+\gamma)}{\gamma}} + \ln\left(r^{-1}\right)^{\frac{2-2\alpha+\gamma}{\alpha}} \right)$$

B Proofs of Section 4

Our goal in this section is to prove Theorem 10 and Proposition 11. To do so, we begin by introducing a lower bound that compares the second moment of ρ with a certain moment from π in order to lower bound the Rényi divergence between ρ and π , thus allowing us to track the evolution of the second moment of the process rather than the Rényi divergence itself.

Lemma 19. Let q > 1, and ρ and π be a pair of measures with $\rho(\|\cdot\|^2) < \infty$ and $\pi(\|\cdot\|^{\frac{2q}{q-1}}) < \infty$, then the *q*-Rényi divergence is lower bounded by

$$R_q(\rho \| \pi) \ge \ln\left(\frac{\rho(\|\cdot\|^2)^{\frac{q}{q-1}}}{\pi(\|\cdot\|^{\frac{2q}{q-1}})}\right).$$

Proof. By [BDK⁺21, Theorem 3.1], we have the following variational representation for the Rényi divergence,

$$R_q(\rho \parallel \pi) \ge \sup_{\phi} \bigg\{ \frac{q}{q-1} \ln \int \exp((q-1)\phi(x)) d\rho(x) - \ln \int \exp(q\phi(x)) d\pi(x) \bigg\},$$

where the supremum runs over all measurable functions. The choice of $\phi(x) = \frac{2}{q-1} \ln(||x||)$ proves the statement of the lemma.

The following lemma uses the gradient bound condition in (4.1) to lower bound the decay rate of the second moment for the Langevin diffusion and LMC.

Lemma 20 (Evolution of the Second Moment). Suppose Eq. (4.1) holds with $\alpha \in [0, 2]$, and $\mathbb{E}||X_0||^2 < \infty$. Then,

$$\frac{\mathrm{d}}{\mathrm{d}t}\mathbb{E}\|X_t\|^2 \ge 2d - 2b\mathbb{E}\Big[\|X_t\|^2\Big]^{\alpha/2}.$$

Similarly, LMC satisfies

$$\mathbb{E}||x_{k+1}||^{2} \ge \mathbb{E}||x_{k}||^{2} - 2bh\mathbb{E}[||x_{k}||^{2}]^{\alpha/2} + 2hd.$$

Proof. We begin by proving the result for LMC. By the independence of ξ_k and x_k , we have

$$\mathbb{E} \|x_{k+1}\|^2 = \mathbb{E} \|x_k - h\nabla V(x_k)\|^2 + 2hd$$

$$\geq \mathbb{E} \|x_k\|^2 - 2h\mathbb{E} \langle \nabla V(x_k), x_k \rangle + 2hd$$

$$\geq \mathbb{E} \|x_k\|^2 - 2bh\mathbb{E} \Big[\|x_k\|^2 \Big]^{\alpha/2} + 2hd,$$

where the last inequality follows from Cauchy-Schwartz and Jensen's inequalities. For the diffusion, it follows from Itô's lemma that

$$||X_t||^2 = -2\int_0^t \langle \nabla V(X_s), X_s \rangle ds + 2td + 2\sqrt{2} \int_0^t \langle X_s, dB_s \rangle.$$
(B.1)

We proceed by showing the last term is a martingale and can be removed once taking expectations. For this, it is sufficient to show that X_t is B_t -integrable or equivalently,

$$\mathbb{E}\Big[\int_0^t \|X_s\|^2 \mathrm{d}s\Big] < \infty. \tag{B.2}$$

Indeed, from Itô's lemma, and Tonelli's theorem,

$$\begin{split} \mathbb{E}\|X_t\|^2 &\leq -2\mathbb{E}\left[\int_0^t \langle \nabla V(X_s), X_s \rangle \mathrm{d}s\right] + 2td + 2\sqrt{2}\mathbb{E}\left[\left(\int_0^t \langle X_s, \mathrm{d}B_s \rangle\right)^2\right]^{1/2} \\ &\leq 2(b+d)t + 2b\mathbb{E}\left[\int_0^t \|X_s\|^2 \mathrm{d}s\right] + 2\sqrt{2}\mathbb{E}\left[\int_0^t \|X_s\|^2 \mathrm{d}s\right]^{1/2} \\ &\leq \sqrt{2} + 2(b+d)t + (2b+\sqrt{2})\int_0^t \mathbb{E}\|X_s\|^2 \mathrm{d}s. \end{split}$$

Define $f(s) \coloneqq \sup_{r \in [0,s]} \mathbb{E}\left[\left\| X_r \right\|^2 \right]$. Then, for any $T \ge t$,

$$f(t) \le \underbrace{\sqrt{2} + 2(b+d)T}_{=:C_1} + \underbrace{(2b+\sqrt{2})}_{=:C_2} \int_0^t f(s) \mathrm{d}s.$$

Then, by Grönwall's lemma,

$$\int_{0}^{t} f(s)ds \le \frac{C_{1}}{C_{2}}(\exp(C_{2}t) - 1) < \infty,$$

and consequently (B.2) is satisfied. Thus, we can remove the last term in (B.1) after taking expectation, and after taking a time derivative obtain

$$\frac{\mathrm{d}}{\mathrm{d}t}\mathbb{E}||X_t||^2 = -2\mathbb{E}[\langle \nabla V(X_t), X_t \rangle] + 2d$$
$$\geq -2b\mathbb{E}||X_t||^{\alpha} + 2d$$
$$\geq -2b\mathbb{E}[||X_t||^2]^{\alpha/2} + 2d,$$

where once again the last inequality is implied by Cuachy-Schwartz and Jensen's inequalities.

Another key ingredient of our proof will be controlling the Rényi or KL divergence using the variance of an isotropic Gaussian initialization, which is provided by the following lemmas.

Lemma 21 (Controlling Rényi Divergence by Initial Variance). Let $\pi \propto e^{-V}$, $Z := \int e^{-V(x)} dx$ and suppose V satisfies (4.1). Then the following holds:

1. If $\alpha = 0$, then for $\sigma^2 \ge (d + \nu)^{-1}$, where we recall $\nu \coloneqq b - d > 0$,

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \parallel \pi) \leq \frac{\nu}{2} \ln \sigma^{2} + \ln \left\{ \left(\frac{Z}{(2\pi)^{d/2}} \right) \left(\frac{d+\nu}{e} \right)^{\frac{d+\nu}{2}} \right\} + \frac{1}{2\sigma^{2}}.$$

2. If $\alpha \in (0,2)$, then for $\sigma^2 \ge b^{-1}$,

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \| \pi) \leq \frac{(b\sigma^{2})^{\frac{2}{2-\alpha}}}{\alpha} + \ln\left(\frac{Z}{(2\pi\sigma^{2})^{d/2}}\right) + \frac{1}{2\sigma^{2}}.$$

Proof. We begin by upper bounding V. Using (4.1), we have that

$$V(x) \le \int_0^1 \|\nabla V(tx)\| \|x\| dt \le \begin{cases} \frac{b}{\alpha} \left((1 + \|x\|^2)^{\alpha/2} - 1 \right) & \text{if } \alpha > 0, \\ \frac{b}{2} \ln(1 + \|x\|^2) & \text{if } \alpha = 0. \end{cases}$$

In the case that $\alpha = 0$, it follows that

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \| \pi) \leq \ln\left(\frac{Z}{(2\pi\sigma^{2})^{d/2}}\right) + \sup_{x} \frac{-\|x\|^{2}}{2\sigma^{2}} + \frac{b}{2}\ln(1+\|x\|^{2}).$$

For $\sigma^2 \geq \frac{1}{a+b}$, the supremum occurs at $||x^*||^2 = b\sigma^2 - 1$, and

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \| \pi) \leq \frac{b-d}{2} \ln \sigma^{2} + \ln\left(\frac{Z}{(2\pi)^{d/2}}\right) + \frac{1}{2\sigma^{2}} + \frac{b}{2} \ln \frac{b}{e}.$$

Now suppose $\alpha > 0$, then

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \| \pi) \leq \ln\left(\frac{Z}{(2\pi\sigma^{2})^{d/2}}\right) + \sup_{x} \frac{-\|x\|^{2}}{2\sigma^{2}} + \frac{b}{\alpha}(1 + \|x\|^{2})^{\alpha/2}.$$

For $\sigma^2 \ge b^{-1}$, the supremum is attained at $||x^*||^2 = (b\sigma^2)^{\frac{2}{2-\alpha}} - 1$, and

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \parallel \pi) \leq \frac{b^{\frac{2}{2-\alpha}}}{\alpha} \sigma^{\frac{2\alpha}{2-\alpha}} + \ln\left(\frac{Z}{(2\pi\sigma^{2})^{d/2}}\right) + \frac{1}{2\sigma^{2}}.$$

We pause here to state the result of Lemma 21 more explicitly for the generalized Cauchy potential V_{ν} and the sublinear potential V_{α} . For the first example, we have the normalizing constant

$$Z_{\nu} \coloneqq \frac{\pi^{d/2} \Gamma(\nu/2)}{\Gamma((d+\nu)/2)} \le \pi^{d/2} \Gamma(\nu/2) \left(\frac{d+\nu}{2e}\right)^{-\frac{d+\nu}{2}+1},$$

where the second inequality holds when $d \ge 2$ using $\Gamma(z) \ge (z/e)^{z-1}$ for $z \ge 1$.

Corollary 22. Consider the measure $\pi_{\nu} \propto \exp(-V_{\nu})$ with $V_{\nu}(x) = \frac{d+\nu}{2}\ln(1+\|x\|^2)$. Then, for $d \geq 2$ and $\sigma^2 \geq (d+\nu)^{-1}$, we have

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \parallel \pi_{\nu}) \leq \frac{\nu}{2} \ln \sigma^{2} + \ln \left(2^{\nu/2} \Gamma(\nu/2) \right) + \ln \left(\frac{d+\nu}{2e} \right)$$

The second example, the sub-linear potential $V_{\alpha} = (1 + ||x||^2)^{\alpha/2} - 1$ with $\alpha \in (0, 1)$, satisfies (4.1) with α , and its normalizing constant can be estimated by

$$Z_{\alpha} \coloneqq \int \exp((1 + \|x\|^2)^{\alpha/2} - 1) dx \le \int \exp(\|x\|^{\alpha}) dx = \frac{\pi^{d/2} d/\alpha \Gamma(d/\alpha)}{\Gamma(d/2 + 1)} \le \pi^{d/2} \left(\frac{d}{\alpha}\right)^{d/\alpha - d/2},$$

where we refer to (C.3) for a proof of the identity, and the second inequality follows from Lemma 32.

Corollary 23. Consider the measure $\pi_{\alpha} \propto \exp(-V_{\alpha})$ with $V_{\alpha} = (1 + ||x||^2)^{\alpha/2}$, $\alpha \in (0, 1)$, and $d \geq 2$. Then for $\sigma^2 \geq 1/\alpha$,

$$R_{\infty}(\mathcal{N}_{\sigma^{2}I_{d}} \parallel \pi_{\alpha}) \leq (\alpha \sigma^{2})^{\frac{\alpha}{2-\alpha}} + \left(\frac{d}{\alpha} - \frac{d}{2}\right) \ln(d/\alpha) - \frac{d}{2} \ln(2\sigma^{2}) + \frac{1}{2\sigma^{2}}.$$

When working with tail growth of order $\alpha = 2$, a large initial variance can lead to infinity Rényi divergence of any order q > 1. Hence, we will instead use the KL divergence as our initial metric in this setting.

Lemma 24 (Controlling KL Divergence by Initial Variance). Let $Z \coloneqq \int e^{-V(x)} dx$ and suppose that V satisfies (4.1) with $\alpha = 2$, then

$$\mathsf{KL}(\mathcal{N}_{\sigma^2 I_d} \parallel \pi) \le \frac{(b\sigma^2 - 1)d}{2} + \ln\left(\frac{Z}{(2\pi\sigma^2)^{d/2}}\right)$$

Proof. First, we upper bound V with

$$V(x) \le \int_0^1 \|\nabla V(tx)\| \|x\| dt \le \frac{b}{2} \|x\|^2.$$

Thus we have

$$\begin{aligned} \mathsf{KL}(\mathcal{N}_{\sigma^2 I_d} \parallel \pi) &= \mathbb{E}_{\mathcal{N}_{\sigma^2 I_d}} \left[\ln \left(\frac{Z}{(2\pi\sigma^2)^{d/2}} \exp\left(\frac{-\|x\|^2}{2\sigma^2} + V(x) \right) \right) \right] \\ &\leq \ln \left(\frac{Z}{(2\pi\sigma^2)^{d/2}} \right) + \frac{(b\sigma^2 - 1)d}{2}. \end{aligned}$$

We are now ready to present the proof of Theorem 10.

Proof. [Proof of Theorem 10] Notice that to lower bound the time or iteration complexity, it suffices to lower bound the time or iteration complexity for one specific initialization with initial divergence less than Δ_0 . Our general strategy will be to use Gaussian initializations with large initial variances, such that Lemma 21 ensures the initial divergence is less than Δ_0 , while the time estimate from Lemma 20 provides the lower bound. Throughout this proof, $Z \coloneqq \int e^{-V}$ denotes the normalizing constant.

1. <u>The case $\alpha = 0$ </u>: Suppose that Δ_0 satisfies

$$\Delta_0 \ge \left\{ 1 + 2\ln\left(\frac{Z((d+\nu)/e)^{\frac{d+\nu}{2}}}{(2\pi)^{d/2}}\right) \right\} \lor \nu \ln\left(\frac{2e\pi(\|\cdot\|^{\frac{2q}{q-1}})^{\frac{q-1}{q}}}{d}\right) \lor \nu.$$
(B.3)

Choose $\rho_0 = \mathcal{N}(0, \sigma^2 I_d)$ with $\sigma^2 = \exp\left(\frac{\Delta_0}{\nu}\right)$. By Lemma 21 we have

$$R_{\infty}(\rho_0 \parallel \pi) \le \frac{\Delta_0}{2} + \ln\left\{ \left(\frac{Z}{(2\pi)^{d/2}}\right) \left(\frac{d+\nu}{e}\right)^{\frac{d+\nu}{2}} \right\} + \frac{1}{2\sigma^2} \le \Delta_0$$

Note that due to Lemma 19, in order to have $R_q(\rho_T \parallel \pi) \leq 1$, T needs to be sufficiently large such that

$$\mathbb{E}\left[\left\|X_{T}\right\|^{2}\right] \leq e^{\frac{q-1}{q}}\pi\left(\left\|\cdot\right\|^{\frac{2q}{q-1}}\right)^{\frac{q-1}{q}}.$$

From Lemma 20 we obtain

$$T \ge \frac{\mathbb{E}\left[\left\|X_{0}\right\|^{2}\right] - \mathbb{E}\left[\left\|X_{T}\right\|^{2}\right]}{2\nu},$$
$$N \ge \frac{\mathbb{E}\left[\left\|X_{0}\right\|^{2}\right] - \mathbb{E}\left[\left\|X_{k}\right\|^{2}\right]}{2h\nu}.$$

and

Consequently,
$$T$$
 needs to satisfy

$$T \ge \frac{d\exp\left(\frac{\Delta_0}{\nu}\right) - e\pi\left(\left\|\cdot\right\|^{\frac{2q}{q-1}}\right)^{\frac{q-1}{q}}}{2\nu} \ge \frac{d\exp\left(\frac{\Delta_0}{\nu}\right)}{4\nu}$$

and N needs to satisfy

$$N \ge \frac{d \exp\left(\frac{\Delta_0}{\nu}\right)}{4h\nu}$$

2. The case $0 < \alpha < 2$: Suppose that Δ_0 satisfies

$$\Delta_0 \ge \left\{ \frac{b^{\frac{\alpha}{2-\alpha}}}{\alpha} \left(1 \vee \frac{(eZ^2)^{1/d}}{2\pi} \right)^{\frac{\alpha}{2-\alpha}} \right\} \vee \left(\frac{2^{\frac{2}{2-\alpha}} eb\pi(\|\cdot\|^{\frac{2q}{q-1}})^{\frac{q-1}{q}}}{\alpha^{2/\alpha-1}d} \right)^{\frac{\alpha}{2-\alpha}} \vee \frac{1}{\alpha}. \tag{B.4}$$

This time, we choose $\rho_0 = \mathcal{N}(0, \sigma^2 I_d)$ with $\sigma^2 = \frac{(\alpha \Delta_0)^{\frac{2-\alpha}{\alpha}}}{b}$. Then, (B.4) ensures $\sigma^2 \ge 1$ and by Lemma 21,

$$R_{\infty}(\rho_0 \parallel \pi) \leq \Delta_0 + \ln\left(\frac{Z}{(2\pi\sigma^2)^{d/2}}\right) + \frac{1}{2\sigma^2} \leq \Delta_0$$
$$\mathbb{E}\left[\|X_0\|^2\right]^{1-\alpha/2} - \mathbb{E}\left[\|X_{\pi}\|^2\right]^{1-\alpha/2}$$

From Lemma 20

$$T \ge \frac{\mathbb{E}\left[\|X_0\|^2\right]^{1-\alpha/2} - \mathbb{E}\left[\|X_T\|^2\right]^{1-\alpha/2}}{b(2-\alpha)}$$

hence we can write

$$T \ge \frac{\left(\frac{\alpha^{2/\alpha-1}}{b}\right)^{1-\alpha/2} d^{1-\alpha/2} \Delta_0^{\frac{(2-\alpha)^2}{2\alpha}} - (e\pi(\|\cdot\|^{\frac{2q}{q-1}})^{\frac{q-1}{q}})^{1-\alpha/2}}{b(2-\alpha)}$$
$$\ge \frac{\left(\frac{\alpha^{2/\alpha-1}}{b}\right)^{1-\alpha/2} d^{1-\alpha/2} \Delta_0^{\frac{(2-\alpha)^2}{2\alpha}}}{2(2-\alpha)b}.$$

For the discrete-time case, from Lemma 20 we obtain

$$\mathbb{E}\left[\left\|x_{k+1}\right\|^{2}\right] \geq \mathbb{E}\left[\left\|x_{k}\right\|^{2}\right] - 2hb\mathbb{E}\left[\left\|x_{k}\right\|^{2}\right]^{\alpha/2}.$$

Let $r_k := \mathbb{E}\left[\|x_k\|^2 \right]$. Suppose $r_k \ge r_{k+1}$, rearranging the inequality above, we obtain

$$2hb \ge r_k^{1-\alpha/2} - r_{k+1}r_k^{-\alpha/2} \ge r_k^{1-\alpha/2} - r_{k+1}^{1-\alpha/2}.$$

On the other hand, when $r_k < r_{k+1}$,

$$2hb > 0 > r_k^{1-\alpha/2} - r_{k+1}^{1-\alpha/2}.$$

Thus the bound holds in either case, and by iterating it we have

$$N \ge \frac{\mathbb{E}\Big[\|X_0\|^2 \Big]^{1-\alpha/2} - \mathbb{E}\Big[\|X_T\|^2 \Big]^{1-\alpha/2}}{hb(2-\alpha)} \ge \frac{\left(\frac{\alpha^{2/\alpha-1}}{\tilde{a}}\right)^{1-\alpha/2} d^{1-\alpha/2} \Delta_0^{\frac{(2-\alpha)^2}{2\alpha}}}{2(2-\alpha)hb}$$

where the second inequality follows analogously to the continuous-time case.

3. <u>The case $\alpha = 2$ </u>: Suppose that Δ_0 satisfies

$$\Delta_0 \ge \frac{bZ^{2/d}e^{2a/d-1}}{4\pi} \lor b\left(e\pi(\|\cdot\|^{\frac{2q}{q-1}})\right)^{\frac{(1+c)(q-1)}{q}}.$$
(B.5)

for any absolute constant c > 0. Choose $\rho_0 = \mathcal{N}(0, \sigma^2 I_d)$ with $\sigma^2 = \frac{2\Delta_0}{bd}$. Then, by Lemma 24 we have

$$\mathsf{KL}(\rho_0 \parallel \pi) \le \Delta_0 - \frac{d}{2} + \ln\left(\frac{Z}{(2\pi\sigma^2)^{d/2}}\right) \le \Delta_0$$

Moreover, from Lemma 20, we have

$$T \ge \frac{\ln\left(\mathbb{E}\left[\|X_0\|^2\right]\right) - \ln\left(\mathbb{E}\left[\|X_T\|^2\right]\right)}{2b} \ge \frac{c\ln\left(\frac{\Delta_0}{b}\right)}{2(1+c)b}$$

Similarly for LMC, when $h < b^{-1}$, we have

$$N \ge \frac{\ln\left(\mathbb{E}\left[\|x_0\|^2\right]\right) - \ln\left(\mathbb{E}\left[\|x_N\|^2\right]\right)}{2hb} \ge \frac{c\ln\left(\frac{\Delta_0}{b}\right)}{2(1+c)b}$$

which completes the proof of the theorem.

In order to prove Proposition 11, we need a sharper control on the decay of the second moment of LMC that does not ignore terms of order $O(h^2)$. In the following lemma, we achieve this control in the radially symmetric setting.

Lemma 25 (A Sharper Evolution Inequality for LMC). Suppose $\mathbb{E}[||x_0||^2] < \infty$, the potential is radially symmetric with $V(x) = f(||x||^2)$ and the function $g : \mathbb{R}_+ \to \mathbb{R}_+$ given by $g(r) = (1 - 2hf'(r))^2 r$ is convex, then for each $k \ge 0$,

$$\mathbb{E}||x_{k+1}||^2 \ge g(\mathbb{E}\left[||x_k||^2\right]) + 2hd.$$

Furthermore, if g is non-decreasing then, $\mathbb{E}||x_k||^2 \ge ||y_k||^2$ for each $k \ge 0$, where we define y_k by

$$y_{k+1} = y_k - \eta \nabla V(y_k), \quad ||y_0||^2 = \mathbb{E}\Big[||x_0||^2\Big].$$

Proof. Using the independence of the Gaussian perturbations and the fact that $\nabla V(x) = 2f'(||x||^2)x$,

$$\mathbb{E}\Big[\|x_{k+1}\|^2\Big] = \mathbb{E}\|x_k - h\nabla V(x_k)\|^2 + 2hd = \mathbb{E}g(\|x_k\|^2) + 2hd.$$

Using the convexity of g along with Jensen's inequality, we conclude that

$$\mathbb{E}\left[\left\|x_{k+1}\right\|^{2}\right] \geq g(\mathbb{E}\left[\left\|x_{k}\right\|^{2}\right]) + 2hd$$

If g is non-decreasing, it follows by comparison that $\mathbb{E}\left[\|x_k\|^2\right] \ge \|y_k\|^2$.

Finally, we can present the proof of Proposition 11.

Proof. [Proof of Proposition 11] First, notice that by Lemma 19, $\inf_k R_q(\mu_k \parallel \pi) < \varepsilon$ is equivalent to $\inf_{k \in \mathbb{N}} \mathbb{E}[\|x_k\|^2] < \sigma_{\varepsilon}^2$. Let z_k be the process defined by the update

$$z_{k+1} = g(z_k) + 2hd, \quad z_0 = \mathbb{E}\left[\|x_0\|^2 \right]$$

so that, using Lemma 25, we have $\mathbb{E}[\|x_k\|^2] \geq z_k$. Thus, if it holds that $\inf_{k\in\mathbb{N}}\mathbb{E}[\|x_k\|^2] < \sigma_{\varepsilon}^2$ then $\inf_{k\in\mathbb{N}} z_k < \sigma_{\varepsilon}^2$ must hold also. If this holds, there must be some $k\in\mathbb{N}$ such that $z_{k+1} \leq \sigma_{\varepsilon}^2$ and $z_k \geq \sigma_{\varepsilon}^2$. Thus, by the fact that g is non-decreasing,

$$g(\sigma_{\varepsilon}^2) + 2hd \le g(z_k) + 2hd = z_{k+1} \le \sigma_{\varepsilon}^2.$$

Rearranging this leads to the bound given in the statement.

C Auxiliary Lemmas

In this section, we prove various moment and tail bounds for generalized Cauchy measures and measures with sublinear potentials, which we use in other proofs of the paper.

Lemma 26. Consider the measure $\tilde{\pi}_{\alpha}(x) \propto \exp(-\lambda \|x\|^{\alpha})$ for $0 < \alpha \leq d$ and $\lambda > 0$. Then, for any p > 0

$$\mathbb{E}_{\tilde{\pi}_{\alpha}}[\|x\|^{p}] = \lambda^{-p/\alpha} \frac{\Gamma\left(\frac{d+p}{\alpha}\right)}{\Gamma\left(\frac{d}{\alpha}\right)} \le \lambda^{-p/\alpha} \left(\frac{d+p}{\alpha}\right)^{\frac{p}{\alpha}}.$$
 (C.1)

Moreover, for $\alpha \in (0,1)$ and $\pi_{\alpha}(x) \propto \exp(-(1+\lambda^{2/\alpha}||x||^2)^{\alpha/2})$,

$$\mathbb{E}_{\pi_{\alpha}}[\|x\|^{p}] \leq \lambda^{-p/\alpha} \frac{e\Gamma\left(\frac{d+p}{\alpha}\right)}{\Gamma\left(\frac{d}{\alpha}\right)}.$$
(C.2)

Proof. First, consider the case of $\lambda = 1$. We begin by computing the normalizing factor of $\tilde{\pi}_{\alpha}$. Using the polar coordinates, we have

$$Z := \int \exp(-\|x\|^{\alpha}) dx = d\omega_d \int_0^\infty \exp(-r^{\alpha}) r^{d-1} dr$$
$$= \frac{d\omega_d}{\alpha} \int_0^\infty \exp(-u) u^{d/\alpha - 1} du$$
$$= \frac{d\omega_d}{\alpha} \Gamma\left(\frac{d}{\alpha}\right).$$
(C.3)

Via similar calculations, we obtain

$$\int \exp(-\|x\|^{\alpha}) \|x\|^{p} \mathrm{d}x = \frac{d\omega_{d}}{\alpha} \Gamma\left(\frac{d+p}{\alpha}\right).$$

Applying Lemma 32 yields,

$$\mathbb{E}_{\tilde{\pi}_{\alpha}}[\|x\|^{p}] = \frac{\Gamma\left(\frac{d+p}{\alpha}\right)}{\Gamma\left(\frac{d}{\alpha}\right)} \le \left(\frac{d+p}{\alpha}\right)^{\frac{p}{\alpha}}$$

Finally, we observe that for $\pi_{\alpha}(x) \propto (1 + ||x||^2)^{\alpha/2}$,

$$\mathbb{E}_{\pi_{\alpha}}[\|x\|^{p}] = \frac{\int \exp\left(-(1+\|x\|^{2})^{\frac{\alpha}{2}}\right) \|x\|^{p} dx}{\int \exp\left(-(1+\|x\|^{2})^{\frac{\alpha}{2}}\right) dx} \le \frac{\int \exp(-\|x\|^{\alpha}) \|x\|^{p} dx}{e^{-1} \int \exp(-\|x\|^{\alpha}) dx} = \frac{e\Gamma\left(\frac{d+p}{\alpha}\right)}{\Gamma\left(\frac{d}{\alpha}\right)}.$$

For the case of $\lambda > 0$, we use the change of variables formula to show that scaling by $\lambda^{-1/\alpha}$ recovers a random variable with the density given by the case with $\lambda = 1$.

Lemma 27. Consider the measure $\pi_{\nu}(x) \propto (1 + ||x||^2)^{-(d+\nu)/2}$ for $\alpha > 0$ with $\nu > p \ge 0$. Then,

$$\mathbb{E}_{\pi_{\nu}}[\|x\|^{p}] = \frac{d}{d+p} \frac{\Gamma\left(\frac{\nu-p}{2}\right)}{\Gamma\left(\frac{\nu}{2}\right)} \frac{\Gamma\left(\frac{d+2+p}{2}\right)}{\Gamma\left(\frac{d+2}{2}\right)} \le \frac{\Gamma\left(\frac{\nu-p}{2}\right)}{\Gamma\left(\frac{\nu}{2}\right)} \left(\frac{d+2+p}{2}\right)^{p/2}.$$

Proof. Recall that the normalizing constant of this measure is given by

$$Z_{d,\nu} \coloneqq \frac{\Gamma(\frac{\nu}{2})\pi^{d/2}}{\Gamma(\frac{\nu+d}{2})} = \frac{\Gamma(\frac{\nu}{2})\Gamma(\frac{d+2}{2})\omega_d}{\Gamma(\frac{\nu+d}{2})}.$$

On the other hand, using the polar coordinates, one can observe

$$Z_{d,\nu} = d\omega_d \int (1+r^2)^{-(\nu+d)/2} r^{d-1} \mathrm{d}r.$$
 (C.4)

We proceed to compute the following

$$\mathbb{E}_{\pi_{\nu}}[\|x\|^{p}] = \frac{d\omega_{d}}{Z_{d,\nu}} \int (1+r^{2})^{-(d+\nu)/2} r^{d+p-1} \mathrm{d}r = \frac{d\omega_{d}}{Z_{d,\nu}} \frac{Z_{d+p,\nu-p}}{(d+p)\omega_{d+p}},$$

where the second equality follows from a change of variables in (C.4). The statement of the lemma follows by an application of Lemma 32.

Lemma 28. The measure
$$\pi_{\alpha} \propto \exp\left(-\left(1+\|x\|^2\right)^{\alpha/2}\right)$$
 with $\alpha \in (0,1)$ satisfies
 $\pi_{\alpha}(\|x\| \ge R) \le e^{\frac{1}{2}}2^{d/\alpha}\exp\left(-\frac{1}{2}\left(1+R^2\right)^{\alpha/2}\right).$ (C.5)

Proof. Using the Markov inequality,

$$\pi_{\alpha}(\|x\| \ge R) = \pi_{\alpha} \left\{ \exp\left(\frac{1}{2} \left(1 + \|x\|^{2}\right)^{\alpha/2}\right) \ge \exp\left(\frac{1}{2} \left(1 + R^{2}\right)^{\alpha/2}\right) \right\}$$
$$\le \exp\left(-\frac{1}{2} \left(1 + R^{2}\right)^{\alpha/2}\right) \mathbb{E}_{\pi_{\alpha}} \left[\exp\left(\frac{1}{2} \left(1 + \|x\|^{2}\right)^{\alpha/2}\right)\right].$$

Using polar coordinates and change of coordinates,

$$\begin{split} \mathbb{E}_{\pi_{\alpha}} \bigg[\exp\bigg(\frac{1}{2} \Big(1 + \|x\|^2\Big)^{\alpha/2} \Big) \bigg] &= \frac{\int \exp\bigg(-\frac{1}{2} \Big(1 + \|x\|^2\Big)^{\alpha/2} \Big) \mathrm{d}x}{\int \exp\bigg(-\Big(1 + \|x\|^2\Big)^{\alpha/2} \Big) \mathrm{d}x} \\ &= \frac{\int_0^\infty r^{d-1} \exp\bigg(-\frac{1}{2} (1 + r^2)^{\alpha/2} \Big) \mathrm{d}r}{\int_0^\infty r^{d-1} \exp\bigg(-(1 + r^2)^{\alpha/2} \Big) \mathrm{d}r} \\ &= \frac{e^{-\frac{1}{2}} \int_0^\infty (u+1)^{2/\alpha-1} \Big[(u+1)^{2/\alpha} - 1\Big]^{(d-2)/2} \exp(-u/2) \mathrm{d}u}{e^{-1} \int_0^\infty (u+1)^{2/\alpha-1} \Big[(u+1)^{2/\alpha} - 1\Big]^{(d-2)/2} \exp(-u) \mathrm{d}u} \\ &\leq e^{\frac{1}{2}} 2^{d/\alpha}, \end{split}$$

where we used the change of variables $u = (1 + r^2)^{\alpha/2} - 1$. This completes the proof.

Lemma 29. The measure $\pi_{\nu}(x) \propto (1 + ||x||^2)^{-(d+\nu)/2}$ satisfies

$$\pi_{\nu}(\|x\| \ge R) \le (\nu+d)^{\nu/2} R^{-\nu}.$$
(C.6)

Proof. Using polar coordinates

$$\pi_{\nu}(\|x\| \ge R) = \frac{1}{Z} \int_{\|x\| \ge R} (1 + \|x\|^2)^{-(d+\nu)/2} \mathrm{d}x = \frac{d\omega_d}{Z} \int_R^\infty (1 + r^2)^{-(d+\nu)/2} r^{d-1} \mathrm{d}r \le \frac{d\omega_d R^{-\nu}}{\nu Z}.$$

And $Z = \frac{\Gamma(\nu/2)\pi^{d/2}}{\Gamma((\nu+d)/2)}$. Hence,

$$\pi_{\nu}(\|x\| \ge R) \le \underbrace{\frac{d\Gamma\left(\frac{\nu+d}{2}\right)}{\nu\Gamma\left(\frac{\nu}{2}\right)\Gamma\left(\frac{d+2}{2}\right)}}_{=:A_{\nu,d}} R^{-\nu}.$$

Suppose $\nu < 2$. Then by Equation (3.1) of [LN13], we have $\frac{\Gamma((d+\nu)/2)}{\Gamma((d+2)/2)} \leq (d/2)^{\nu/2-1}$. Moreover, when $\nu \geq 2$, using Lemma 32 we have $\frac{\Gamma((d+\nu)/2)}{\Gamma((d+2)/2)} \leq ((d+\nu)/2)^{\nu/2-1}$. Consequently,

$$A_{\nu,d} = \frac{d^{\nu/2} (1/2)^{\nu/2-1} (1 \vee (1+\nu/d)^{\nu/2-1})}{\nu \Gamma(\nu/2)} \le \frac{(d+\nu)^{\nu/2}}{2^{\nu/2} \Gamma(\nu/2+1)} \le (d+\nu)^{\nu/2},$$

where we used the fact that $2^{\nu/2}\Gamma(\nu/2+1) \ge \Gamma(1) = 1$.

The following Lemma, adapted from [CEL⁺22], shows the existence of isotropic Gaussian initializations such that $R_q(\mu_0 \parallel \pi), R_q(\mu_0 \parallel \hat{\pi}) = \tilde{O}(d)$.

Lemma 30. Let $\pi(x) \propto \exp(-V(x))$ such that ∇V is s-Hölder continuous and $\nabla V(0) = 0$. Define $\hat{\mu}$ as in Proposition 14. Let $\mu_0 = \mathcal{N}(0, (2L+1)^{-1}I_d)$. Then,

$$R_{\infty}(\mu_0 \parallel \pi) \le 2 + L + V(0) - \min_x V(x) + \frac{d}{2} \ln(12m^2 L),$$
(C.7)

$$R_{\infty}(\mu_0 \parallel \hat{\pi}) \le 3 + L + V(0) - \min_x V(x) + \frac{d}{2} \ln(12(m + 6144T)^2 L).$$
(C.8)

Proof. The Lemma is directly based on Lemmas 30 and 31 of [CEL⁺22].

The following lemma translates a bound on $R_2(\mu_0 || \pi)$ to a bound on $R_2(\mu_0 || \hat{\pi})$ when μ_0 is some isotropic Gaussian measure. As we only calculate the former quantity for our model examples, we use this lemma to establish similar bounds for the latter.

Lemma 31. Suppose $\pi \propto \exp(-V(x))$ and $\hat{\pi} \propto \exp(-\hat{V}(x))$ with $\hat{V}(x) = V(x) + \frac{\gamma}{2} \max\{\|x\| - R, 0\}^2$

for some $\gamma, R > 0$. Then, for any $\sigma^2 \leq \frac{1}{\gamma}$ we have

$$R_2(\mathcal{N}_{\sigma^2 I_d} \| \hat{\pi}) \le d \ln 2 + R_2(\mathcal{N}_{2\sigma^2 I_d} \| \pi).$$

Proof. Let $Z \coloneqq \int \exp(-V(x)) dx$ and $\hat{Z} \coloneqq \int \exp(-\hat{V}(x)) dx$. Notice that $V \leq \hat{V}$, thus $\hat{Z} \leq Z$. Therefore,

$$R_{2}(\mathcal{N}_{\sigma^{2}I_{d}} \| \hat{\pi}) = \ln\left(\frac{\hat{Z}}{(2\pi\sigma^{2})^{d}} \int \exp\left(\frac{-\|x\|^{2}}{\sigma^{2}} + V(x) + \frac{\gamma}{2}\max\{\|x\| - R, 0\}^{2}\right) dx\right)$$
$$\leq \ln\left(\frac{Z}{(2\pi\sigma^{2})^{d}} \int \exp\left(\frac{-\|x\|^{2}}{2\sigma^{2}} + V(x)\right) dx\right)$$
$$= d\ln 2 + R_{2}(\mathcal{N}_{2\sigma^{2}I_{d}} \| \pi).$$

Lemma 32 ([LN13, Theorem 3.1]). Suppose $x \ge y \ge 1$, then

$$y^{x-y} \le \frac{\Gamma(x)}{\Gamma(y)} \le x^{x-y}.$$
(C.9)

Lemma 33. Suppose $Z \ge 0$ is a non-negative random variable. Then, for any $p \ge 2$,

$$\mathbb{E}[Z^p] \ge \mathbb{E}[|Z - \mathbb{E}[Z]|^p].$$

Proof. Normalize Z such that $\mathbb{E}[Z] = 1$. Using the inequality $Z^p - |Z - 1|^p \ge Z - 1$ for every $Z \ge 0$ and $p \ge 2$ and taking expectations proves the lemma.