
VarFlow: Variational Distillation Through Score-Rule Matching of Noisy Distributions

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Abstract

Diffusion models achieve remarkable generative performance but are hampered by slow, iterative inference. Model distillation seeks to train a fast student generator. Variational Score Distillation (VSD) offers a principled KL-divergence minimization framework for this task. This method cleverly avoids computing the teacher model’s Jacobian, but its student gradient relies on the score of the student’s own noisy marginal distribution, $\nabla_{\mathbf{x}_t} \log p_{\phi,t}(\mathbf{x}_t)$. VSD thus requires approximations, such as training an auxiliary network to estimate this score. These approximations can introduce biases, cause training instability, or lead to an incomplete match of the target distribution, potentially focusing on conditional means rather than broader distributional features. We introduce **VarFlow**, a novel distillation method based on a framework we term *Score-Rule Variational Distillation (SRVD)* framework. VarFlow trains a one-step generator $g_\phi(\mathbf{z})$ by directly minimizing an energy distance (derived from the strictly proper energy score) between the student’s induced noisy data distribution $p_{\phi,t}(\mathbf{x}_t)$ and the teacher’s target noisy distribution $q_t(\mathbf{x}_t)$. This objective is estimated entirely using samples from these two distributions. Crucially, VarFlow bypasses the need to compute or approximate the intractable student score. By directly matching the full noisy marginal distributions, VarFlow aims for a more comprehensive and robust alignment between student and teacher, offering an efficient and theoretically grounded path to high-fidelity one-step generation.

1 Introduction

Diffusion models [Sohl-Dickstein et al. \[2015\]](#), [Song and Ermon \[2019\]](#), [Ho et al. \[2020\]](#), [Song et al. \[2021\]](#) represent a significant advance in generative modeling. They have demonstrated state-of-the-art capabilities in producing high-fidelity and diverse samples across various domains, including image synthesis [Rombach et al. \[2022\]](#), [Ramesh et al. \[2022\]](#), [Ho et al. \[2022\]](#), [Shao et al. \[2023\]](#), audio generation [Liu et al. \[2023\]](#), and 3D content creation [Poole et al. \[2022\]](#), [Wang et al. \[2024\]](#). These models operate by systematically adding noise to data samples in a forward process and then learning to reverse this process, enabling generation from an initial noise distribution. However, the reverse process is iterative, often requiring tens to thousands of sequential steps. This computational burden makes inference slow and limits their practical use in scenarios demanding real-time generation or operating under resource constraints.

To mitigate this high inference cost, various distillation techniques have been developed [Meng et al. \[2023\]](#), [Kang et al. \[2024\]](#), [Salimans and Ho \[2022\]](#), [Shao et al. \[2025\]](#). The goal of these methods is to compress the generative capabilities of a large, multi-step teacher diffusion model into a smaller, faster student model. Ideally, this student model can generate high-quality samples in significantly fewer steps, often in a single forward pass.

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One prominent family of distillation methods aims to match distributions induced by the student and teacher models at different noise levels. For example, Score Distillation Sampling (SDS) [Poole et al. \[2022\]](#) and its variants guide a student generator $g_\phi(z)$ using a pre-trained teacher diffusion model, typically represented by its noise predictor $\epsilon_{\text{teacher}}$. A direct approach to optimize g_ϕ might involve a denoising score matching (DSM)-like loss: generate x_0 using g_ϕ , noise it to x_t , and penalize differences between $\epsilon_{\text{teacher}}(x_t, t)$ and the actual noise ϵ . However, the gradient of such a loss with respect to the student's parameters ϕ can involve the Jacobian of the teacher model, $\nabla_{x_t} \epsilon_{\text{teacher}}(x_t, t)$. Computing this Jacobian is often prohibitively expensive for large teacher models and can be numerically unstable. SDS employs a stop-gradient approximation to make its gradient tractable, effectively treating $\epsilon_{\text{teacher}}(x_t, t) - \epsilon$ as a gradient direction for x_0 . While practical, this update may not correspond to minimizing a well-defined distributional objective.

Variational Score Distillation (VSD) [Wang et al. \[2024\]](#) provides a more theoretically grounded approach. VSD seeks to minimize the KL divergence $D_{\text{KL}}(p_{\phi,t}(x_t) \parallel q_t(x_t))$ between the student's induced noisy data distribution $p_{\phi,t}(x_t)$ and the teacher's target noisy distribution $q_t(x_t)$ across various noise levels t . The gradient of this KL divergence with respect to the student generator's parameters ϕ , $\nabla_\phi D_{\text{KL}}(p_{\phi,t} \parallel q_t)$, involves the difference between the student's score $\nabla_{x_t} \log p_{\phi,t}(x_t)$ and the teacher's score $\nabla_{x_t} \log q_t(x_t)$ (see Equation (6)). The teacher's score is readily available from the pre-trained model $\epsilon_{\text{teacher}}$. A key advantage of the VSD formulation is that its gradient avoids the problematic teacher Jacobian $\nabla_{x_t} \epsilon_{\text{teacher}}$. However, a new challenge arises: the student's score $s_{\phi,t}(x_t) = \nabla_{x_t} \log p_{\phi,t}(x_t)$ is itself intractable. This intractability stems from $p_{\phi,t}(x_t)$ being a marginal distribution, $p_{\phi,t}(x_t) = \int q_t(x_t | g_\phi(z)) p_z(z) dz$, obtained by integrating over the student's latent input z . Consequently, VSD must approximate this student score. Common strategies include training an auxiliary score network ϵ_{aux} (parameterized, for instance, by ω) via DSM on samples generated by g_ϕ and then noised. Another approach involves simpler approximations, such as replacing the marginal student score $s_{\phi,t}(x_t)$ with the score of the conditional distribution $q_t(x_t | g_\phi(z^*))$ for a specific z^* . While these approximations make VSD practical, they can introduce biases, lead to training instabilities (such as those from alternating optimization of ϕ and ω if ϵ_{aux} is used), or result in an incomplete match of the target distribution, potentially focusing more on aligning conditional means rather than broader distributional characteristics.

To address the limitations associated with approximating the student score, we propose VarFlow, a novel distillation method based on a framework we term Score-Rule Variational Distillation. VarFlow trains a fast, single-step generator $g_\phi(z)$ by directly minimizing a statistical distance between the student's induced noisy data distribution $p_{\phi,t}(x_t)$ and the teacher's target noisy distribution $q_t(x_t)$, as illustrated conceptually in Figure 1. The core innovation of VarFlow is its use of an objective derived from strictly proper scoring rules—specifically, the energy score, which leads to minimizing the energy distance [Gneiting and Raftery \[2007\]](#), [Székely et al. \[2004\]](#). This principled choice allows the distillation objective to be estimated entirely using samples drawn from both the student and teacher noisy distributions. As a result, VarFlow completely bypasses the need to compute or approximate the intractable student score $\nabla_{x_t} \log p_{\phi,t}(x_t)$. By directly matching the full distributions $p_{\phi,t}(x_t)$ and $q_t(x_t)$ via an Integral Probability Metric (IPM) like the energy distance, VarFlow aims for a more comprehensive alignment of the student's behavior with the teacher's, compared to methods that concentrate on score or conditional mean matching. This approach is inspired by the successful application of scoring rules in learning complex conditional distributions, as seen in conceptual Distributional Diffusion Models (DDMs, discussed in Section 3.3 and Section D). VarFlow adapts this principle to the distillation context, presenting an efficient and principled alternative for deriving high-quality, rapid generative models from pre-trained diffusion teachers.

Our main contributions are:

1. **A Novel Distillation Method (VarFlow):** We propose VarFlow, which operationalizes a Score-Rule Variational Distillation (SRVD) framework. VarFlow trains a single-step generator by minimizing the energy distance between the noisy data distributions induced by the student and the teacher. This objective is estimated purely from samples, thereby circumventing the need for student score computation or approximation.
2. **Theoretical Grounding and Robustness:** We provide theoretical analysis establishing the consistency of the VarFlow objective. We argue that its sample-based nature and avoidance of score approximations can lead to more stable and straightforward training dynamics.

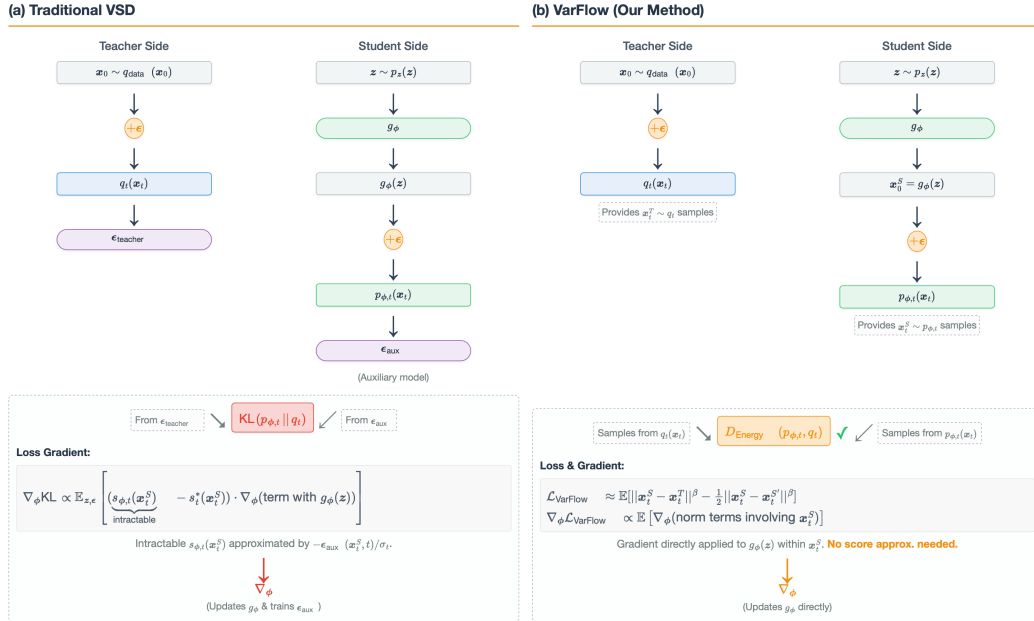


Figure 1: Conceptual comparison. (a) Variational Score Distillation (VSD) minimizes KL divergence. Its gradient requires the student’s marginal score $\nabla_{x_t} \log p_{\phi,t}(x_t)$, which is typically intractable and approximated (e.g., via an auxiliary network ϵ_{aux} trained to predict noise corresponding to this score). (b) VarFlow, our Score-Rule Variational Distillation (SRVD) method, directly minimizes the energy distance between the student’s noisy marginal $p_{\phi,t}(x_t)$ and the teacher’s $q_t(x_t)$. This objective is estimated using only samples from these distributions, bypassing the need for student score estimation and its associated approximations.

3. **Principled Alternative to KL-based VSD:** VarFlow offers a robust alternative to KL-divergence based distillation. By directly optimizing an IPM, it aims to capture fuller distributional characteristics, potentially leading to higher-fidelity student models.

2 Related Work

Our method, VarFlow, distills large diffusion models into one-step generators by directly matching the noisy marginal distributions of a student and teacher using an energy distance objective. This places our work at the intersection of diffusion model distillation, generative models using integral probability metrics, and generative modeling with scoring rules.

2.1 Distillation of Diffusion Models

A primary challenge in diffusion model distillation is removing the slow, iterative sampling process. Variational Score Distillation (VSD) Wang et al. [2024] offers a principled approach by minimizing the KL-divergence between the student and teacher’s noisy marginals. However, VSD requires approximating the intractable student score function, $\nabla_{x_t} \log p_{\phi,t}(x_t)$, which can introduce bias.

Recent works have demonstrated that effective distillation is possible without direct score supervision Zhang et al. [2025]. These methods show that while the teacher’s score function is not essential, initializing the student with the teacher’s weights is crucial for transferring learned feature representations Zhang et al. [2025]. VarFlow shares the goal of avoiding explicit score functions. However, instead of relying on initialization, we propose directly minimizing a well-defined Integral Probability Metric (IPM) between the marginal distributions, offering a more direct path to distributional matching.

2.2 Integral Probability Metrics in Generative Models

Using statistical distances to match distributions is a well-established concept, particularly in Generative Adversarial Networks (GANs). An IPM measures the distance between two probability

distributions, with a prominent example being the Maximum Mean Discrepancy (MMD), which is based on a Reproducing Kernel Hilbert Space (RKHS) [Gretton et al. \[2012\]](#). MMD GAN, for instance, replaces the standard discriminator with a two-sample MMD test, which can stabilize training and mitigate issues like mode collapse by training the generator to minimize the MMD between generated and real samples [Li et al. \[2017\]](#).

VarFlow applies this philosophy to diffusion distillation. We replace the KL-divergence objective found in VSD with an energy distance objective—which is also an IPM—thereby inheriting the benefits of direct distributional matching in a distillation context.

2.3 Scoring Rules in Diffusion Models

Our use of the energy distance is inspired by the application of scoring rules for learning *conditional* distributions in multi-step diffusion models. For example, Distributional Diffusion Models (DDM) [De Bortoli et al. \[2025\]](#) use proper scoring rules to learn the full conditional posterior $q(x_0|x_t)$, not just its mean. This allows for more diverse sampling at each reverse step and can accelerate inference by enabling larger step sizes [De Bortoli et al. \[2025\]](#).

However, the distinction is critical: DDM improves a **multi-step** sampler by enhancing each conditional step, whereas VarFlow is designed for **one-step** generation. Our method matches the **noisy marginal distributions**, $p_{\phi,t}(x_t)$ and $q_t(x_t)$, which allows the objective to be estimated purely from samples and bypasses the student score problem. While methods like DDM [De Bortoli et al. \[2025\]](#) and Reverse Markov Learning [Shen et al. \[2025\]](#) also use scoring rules, they remain multi-step frameworks focused on conditional distributions, unlike VarFlow’s one-step, marginal-matching approach. Other works like Inductive Moment Matching (IMM) [Sun et al. \[2025\]](#) also validate distributional matching for efficient generation, a principle VarFlow adapts for distillation.

The following table summarizes the key distinctions:

Method	Learns	Distribution Matched	Inference Process
RML	Sequence of generators $\{g_t(x_t, \epsilon)\}_{t=0}^T$	Reverse Conditional: $p(x_{t-1} x_t)$	T-step sequential
DDM	Conditional posterior generator $G_\theta(x_t, t, \xi)$	Conditional Posterior: $p(x_0 x_t)$	Multi-step (DDIM-like)
VarFlow (Ours)	Single one-step generator $g_\phi(z)$	Noisy Marginal: $p_{\phi,t}(x_t)$	One-step parallel

This distinction is crucial. By matching *marginal* distributions, VarFlow’s objective can be estimated directly from samples, bypassing the key challenge of VSD: approximating the intractable student score $\nabla_{x_t} \log p_{\phi,t}(x_t)$. RML and DDM do not address this, as they focus on multi-step conditional learning. Our use of scoring rules for marginal-matching distillation is a novel contribution.

3 Preliminaries

We review key concepts: diffusion models, proper scoring rules, the idea of distributional posterior learning as a conceptual precursor, and Variational Score Distillation. A summary of notation is in Table 5 (Section B).

3.1 Diffusion Models and Score Matching

Diffusion models [Ho et al. \[2020\]](#), [Song et al. \[2021\]](#) define a forward noising process that gradually transforms a data sample $x_0 \sim q_{\text{data}}(x_0)$ into approximate Gaussian noise over $t \in [0, T]$. A common formulation, based on the Variance Preserving (VP) SDE [Song et al. \[2021\]](#), defines the conditional distribution of x_t given x_0 as:

$$q_t(x_t | x_0) = \mathcal{N}(x_t; \sqrt{\bar{\alpha}_t} x_0, \sigma_t^2 \mathbf{I}), \quad (1)$$

where $\bar{\alpha}_t$ decreases with t and $\sigma_t^2 = 1 - \bar{\alpha}_t$. Generative modeling involves learning to reverse this process, typically by estimating the score function $\nabla_{x_t} \log q_t(x_t)$, where $q_t(x_t) = \mathbb{E}_{x_0} [q_t(x_t|x_0)]$.

In practice, a network $\epsilon_\theta(\mathbf{x}_t, t)$ predicts the noise ϵ from $\mathbf{x}_t = \sqrt{\alpha_t}\mathbf{x}_0 + \sigma_t\epsilon$, often via Denoising Score Matching (DSM):

$$\mathcal{L}_{\text{DSM}}(\theta) = \mathbb{E}_{t, \mathbf{x}_0, \epsilon} [w(t) \|\epsilon_\theta(\sqrt{\alpha_t}\mathbf{x}_0 + \sigma_t\epsilon, t) - \epsilon\|_2^2]. \quad (2)$$

The learned ϵ_θ relates to the score s_θ via $s_\theta(\mathbf{x}_t, t) \approx -\epsilon_\theta(\mathbf{x}_t, t)/\sigma_t$. Details are in Section C.1.

3.2 Proper Scoring Rules

Scoring rules evaluate probabilistic forecasts [Gneiting and Raftery \[2007\]](#). A scoring rule $\mathcal{S}(P, y)$ assigns a score when forecast P is issued and outcome y occurs. The expected score for $Y \sim Q$ is $\mathcal{S}(P, Q) = \mathbb{E}_{Y \sim Q}[\mathcal{S}(P, Y)]$.

Definition 3.1 (Strictly Proper Scoring Rule). A scoring rule $\mathcal{S}$ is *strictly proper* if $\mathcal{S}(Q, Q) \geq \mathcal{S}(P, Q)$ for all P, Q , with equality iff $P = Q$ (a.e.).

Strict propriety ensures forecasters report true belief. Minimizing $-\mathcal{S}(P_\lambda, Q)$ drives P_λ towards Q .

Energy Score and Energy Distance. A key example is the *energy score* [Gneiting and Raftery \[2007\]](#), [Székely et al. \[2004\]](#). For $\beta \in (0, 2)$:

$$\mathcal{S}_{\text{Energy}}^{(\beta)}(P, y) = \frac{1}{2} \mathbb{E}_{X, X' \stackrel{\text{iid}}{\sim} P} [\|X - X'\|^\beta] - \mathbb{E}_{X \sim P} [\|X - y\|^\beta]. \quad (3)$$

This rule is strictly proper for distributions P with finite β -th moments. Optimizing a model P to match a target Q by minimizing the negative expected energy score, $-\mathbb{E}_{Y \sim Q}[\mathcal{S}_{\text{Energy}}^{(\beta)}(P, Y)]$, is equivalent (up to terms constant in P) to minimizing the *energy distance* $D_{\text{Energy}}^{(\beta)}(P, Q)^2$. The energy distance, detailed in Section C.3 along with the explicit form of the negative expected score (Equation (17)), is an Integral Probability Metric (IPM) defined as:

$$D_{\text{Energy}}^{(\beta)}(P, Q)^2 = 2\mathbb{E}_{X \sim P, Y \sim Q} \|X - Y\|^\beta - \mathbb{E}_{X, X' \stackrel{\text{iid}}{\sim} P} \|X - X'\|^\beta - \mathbb{E}_{Y, Y' \stackrel{\text{iid}}{\sim} Q} \|Y - Y'\|^\beta. \quad (4)$$

It defines a metric on distributions with finite β -moments.

3.3 Conceptual Inspiration: Distributional Learning for Posterior Estimation (DDM)

Standard diffusion models often predict the conditional mean $\mathbb{E}[\mathbf{x}_0|\mathbf{x}_t]$. However, the true posterior $q(\mathbf{x}_0|\mathbf{x}_t)$ can be richer. Distributional Diffusion Models (DDMs, see Section D) aim to learn this entire posterior $p_\theta(\cdot|\mathbf{x}_t, t)$ using a conditional generator $G_\theta(\mathbf{x}_t, t, \xi)$, where ξ is auxiliary noise. The distribution of $G_\theta(\mathbf{x}_t, t, \xi)$ over ξ defines $p_\theta(\cdot|\mathbf{x}_t, t)$. Strictly proper scoring rules, like the energy score, can train G_θ by minimizing an expected negative score $-\mathbb{E}_{(\mathbf{x}_0, \mathbf{x}_t)}[\mathcal{S}(p_\theta(\cdot|\mathbf{x}_t, t)(\cdot|\mathbf{x}_t, t), \mathbf{x}_0)]$. Crucially, the gradient of this objective with respect to θ does *not* require computing the score $\nabla_{\mathbf{x}_0} \log p_\theta(\cdot|\mathbf{x}_t, t)(\hat{\mathbf{x}}_0|\mathbf{x}_t, t)$. This score-free distributional matching is a key inspiration for VarFlow.

3.4 Variational Score Distillation (VSD)

VSD [Wang et al. \[2024\]](#) distills a teacher diffusion model into a fast generator $g_\phi(\mathbf{z})$, where $\mathbf{z} \sim p_z(\mathbf{z})$. The student g_ϕ induces a noisy distribution $p_{\phi, t}(\mathbf{x}_t) = \mathbb{E}_{\mathbf{z} \sim p_z(\mathbf{z})}[\mathcal{N}(\mathbf{x}_t; \sqrt{\alpha_t}g_\phi(\mathbf{z}), \sigma_t^2\mathbf{I})]$. VSD trains g_ϕ by minimizing $D_{\text{KL}}(p_{\phi, t}(\mathbf{x}_t) || q_t(\mathbf{x}_t))$ averaged over time t :

$$\mathcal{L}_{\text{VSD-KL}}(\phi) = \mathbb{E}_{t \sim \text{Unif}[0, T]} [\tilde{w}(t) D_{\text{KL}}(p_{\phi, t}(\mathbf{x}_t) || q_t(\mathbf{x}_t))]. \quad (5)$$

The gradient (see Section C.2 or [Wang et al. \[2024\]](#)) is:

$$\nabla_\phi \mathcal{L}_{\text{VSD-KL}} = \mathbb{E}_{t, \mathbf{z}, \epsilon} [\tilde{w}(t) (s_{\phi, t}(\mathbf{x}_t) - s_t^*(\mathbf{x}_t)) \cdot \nabla_\phi (\sqrt{\alpha_t}g_\phi(\mathbf{z}))], \quad (6)$$

where $\mathbf{x}_t = \sqrt{\alpha_t}g_\phi(\mathbf{z}) + \sigma_t\epsilon$, and $s_t^*(\mathbf{x}_t) \approx -\epsilon_{\text{teacher}}(\mathbf{x}_t, t)/\sigma_t$. The challenge is the intractable student score $s_{\phi, t}(\mathbf{x}_t) = \nabla_{\mathbf{x}_t} \log p_{\phi, t}(\mathbf{x}_t)$. VSD approximates it, e.g., by training an auxiliary network $\epsilon_{\text{aux}}(\mathbf{x}_t, t)$ or using simpler conditional score approximations. These can be sources of issues, which VarFlow avoids.

4 VarFlow: Score-Rule Variational Distillation

We now introduce VarFlow, our method for distilling a pre-trained teacher diffusion model into a fast, single-step generator $g_\phi(z)$. VarFlow operationalizes a Score-Rule Variational Distillation (SRVD) framework. It leverages proper scoring rules to directly minimize a statistical distance between the teacher's and student's noisy data distributions, crucially avoiding the need for student score function estimation or approximation.

4.1 Motivation and Setup

The objective is to train $g_\phi(z)$, where $z \sim p_z(z)$, such that its output $\mathbf{x}_0^S = g_\phi(z)$ resembles data from $q_{\text{data}}(\mathbf{x}_0)$. This student induces a noisy distribution $p_{\phi,t}(\mathbf{x}_t)$:

$$p_{\phi,t}(\mathbf{x}_t|z) = \mathcal{N}(\mathbf{x}_t; \sqrt{\alpha_t}g_\phi(z), \sigma_t^2\mathbf{I}), \quad \text{so} \quad p_{\phi,t}(\mathbf{x}_t) = \mathbb{E}_{z \sim p_z(z)}[p_{\phi,t}(\mathbf{x}_t|z)]. \quad (7)$$

The teacher implies a target noisy distribution $q_t(\mathbf{x}_t)$ from real data $\mathbf{x}_0 \sim q_{\text{data}}(\mathbf{x}_0)$:

$$q_t(\mathbf{x}_t|\mathbf{x}_0) = \mathcal{N}(\mathbf{x}_t; \sqrt{\alpha_t}\mathbf{x}_0, \sigma_t^2\mathbf{I}), \quad \text{so} \quad q_t(\mathbf{x}_t) = \mathbb{E}_{\mathbf{x}_0 \sim q_{\text{data}}(\mathbf{x}_0)}[q_t(\mathbf{x}_t|\mathbf{x}_0)]. \quad (8)$$

VarFlow proposes minimizing the energy distance between $p_{\phi,t}(\mathbf{x}_t)$ and $q_t(\mathbf{x}_t)$, estimated using only samples.

4.2 VarFlow Objective using Energy Score

VarFlow trains g_ϕ by minimizing the energy distance (Equation (4)) between $p_{\phi,t}(\mathbf{x}_t)$ and $q_t(\mathbf{x}_t)$, averaged over t . The *VarFlow Energy Distillation Loss* uses the negative expected energy score (Section C.3, Equation (17)). Minimizing this is equivalent to minimizing energy distance:

$$\mathcal{L}_{\text{VarFlow}}^{(\beta)}(\phi) = \mathbb{E}_{t \sim \text{Unif}[0,T]} \left[\tilde{w}(t) \left(\mathbb{E}_{\substack{\mathbf{x}_t^S \sim p_{\phi,t} \\ \mathbf{x}_t^T \sim q_t}} [\|\mathbf{x}_t^S - \mathbf{x}_t^T\|^\beta] - \frac{1}{2} \mathbb{E}_{\substack{\mathbf{x}_t^S \sim p_{\phi,t} \\ \mathbf{x}_t^{S'} \sim p_{\phi,t}}} [\|\mathbf{x}_t^S - \mathbf{x}_t^{S'}\|^\beta] \right) \right].$$

Here, $\beta \in (0, 2)$ and $\tilde{w}(t)$ is a time weighting. The term $-\frac{1}{2} \mathbb{E}_{\mathbf{x}_t^T, \mathbf{x}_t^{T'} \sim q_t} [\|\mathbf{x}_t^T - \mathbf{x}_t^{T'}\|^\beta]$ from full energy distance (Equation (4)) is omitted as it's constant w.r.t. ϕ .

Monte Carlo Estimation. For a sampled t : 1. Sample $\{z^{(k)}\}_{k=1}^K \stackrel{\text{iid}}{\sim} p_z$; $\mathbf{x}_0^{S,(k)} = g_\phi(z^{(k)})$. 2. Sample $\{\epsilon^{S,(k)}\}_{k=1}^K \stackrel{\text{iid}}{\sim} \mathcal{N}(\mathbf{0}, \mathbf{I})$; $\mathbf{x}_t^{S,(k)} = \sqrt{\alpha_t}\mathbf{x}_0^{S,(k)} + \sigma_t\epsilon^{S,(k)}$. 3. Sample $\{\mathbf{x}_0^{T,(k)}\}_{k=1}^K \stackrel{\text{iid}}{\sim} q_{\text{data}}$. 4. Sample $\{\epsilon^{T,(k)}\}_{k=1}^K \stackrel{\text{iid}}{\sim} \mathcal{N}(\mathbf{0}, \mathbf{I})$; $\mathbf{x}_t^{T,(k)} = \sqrt{\alpha_t}\mathbf{x}_0^{T,(k)} + \sigma_t\epsilon^{T,(k)}$. An unbiased U-statistic estimate for the bracketed term in Equation (8) (for $K \geq 2$):

$$\hat{L}_t(\phi) = \tilde{w}(t) \left(\frac{1}{K^2} \sum_{i=1}^K \sum_{j=1}^K \|\mathbf{x}_t^{S,(i)} - \mathbf{x}_t^{T,(j)}\|^\beta - \frac{1}{K(K-1)} \sum_{i=1}^K \sum_{\substack{j=1 \\ j \neq i}}^K \frac{1}{2} \|\mathbf{x}_t^{S,(i)} - \mathbf{x}_t^{S,(j)}\|^\beta \right). \quad (9)$$

A simpler (paired) estimator for the cross-term, which is often used but may have higher variance:

$$\hat{L}_t^{\text{paired}}(\phi) = \tilde{w}(t) \left(\frac{1}{K} \sum_{k=1}^K \|\mathbf{x}_t^{S,(k)} - \mathbf{x}_t^{T,(k)}\|^\beta - \frac{1}{K(K-1)} \sum_{k=1}^K \sum_{\substack{j=1 \\ j \neq k}}^K \frac{1}{2} \|\mathbf{x}_t^{S,(k)} - \mathbf{x}_t^{S,(j)}\|^\beta \right). \quad (10)$$

The gradient $\nabla_\phi \hat{L}_t(\phi)$ (or $\nabla_\phi \hat{L}_t^{\text{paired}}(\phi)$) is computed via backpropagation through g_ϕ as it appears in $\mathbf{x}_t^{S,(i)}$ and $\mathbf{x}_t^{S,(j)}$.

4.3 Advantages and Connections

The VarFlow method, through its SRVD framework, offers several key advantages:

- **No Explicit Student Score Computation:** A significant advantage is that VarFlow avoids the computational and approximation challenges associated with the student score in VSD. Traditional VSD relies on the gradient (Equation (6)), which includes the term $s_{\phi,t}(\mathbf{x}_t) = \nabla_{\mathbf{x}_t} \log p_{\phi,t}(\mathbf{x}_t)$. This student score is the gradient of the log-density of the student's noisy marginal distribution $p_{\phi,t}(\mathbf{x}_t) = \int p_{\phi,t}(\mathbf{x}_t | g_{\phi}(\mathbf{z})) p_{\mathbf{z}}(\mathbf{z}) d\mathbf{z}$.
- **Direct Distribution Matching via IPMs:** VarFlow minimizes a well-defined statistical distance (the energy distance) between the full distributions $p_{\phi,t}(\mathbf{x}_t)$ and $q_t(\mathbf{x}_t)$. This promotes a comprehensive alignment of their characteristics.
- **Simplicity and Potential for Enhanced Stability:** Avoiding explicit score terms and their approximations can lead to simpler implementation and potentially more stable training.
- **Flexibility with Teacher Information:** VarFlow uses samples from $q_t(\mathbf{x}_t)$ (noised real data) and does not require explicit access to the teacher's score $\epsilon_{\text{teacher}}$ during g_{ϕ} optimization.

Theorem 4.1 (Consistency of VarFlow Objective). *Assume $\tilde{w}(t) > 0$ a.e. on $[0, T]$ and t is sampled with full support. Let $p_{\phi,t}$ be the student's noisy marginal and q_t be the teacher's. If using energy score with $\beta \in (0, 2)$ in $\mathcal{L}_{\text{VarFlow}}^{(\beta)}(\phi)$ (Equation (8)), the loss is minimized iff $p_{\phi,t}(\mathbf{x}_t) = q_t(\mathbf{x}_t)$ for a.e. $t \in [0, T]$ (assuming finite β -moments and sufficient capacity for g_{ϕ}).*

Proof. A sketch is provided here, full details in Section E.1. The VarFlow loss is a weighted integral of terms $-\tilde{w}(t) \mathcal{S}_{\text{Energy}}^{(\beta)}(p_{\phi,t}, q_t)$. Minimizing this w.r.t. $p_{\phi,t}$ is equivalent to minimizing $\frac{1}{2} \tilde{w}(t) D_{\text{Energy}}^{(\beta)}(p_{\phi,t}, q_t)^2$ (plus constant terms). Since $D_{\text{Energy}}^{(\beta)}$ is a metric for $\beta \in (0, 2)$, it's non-negative and zero iff distributions match. Given $\tilde{w}(t) > 0$, overall loss is minimized iff $D_{\text{Energy}}^{(\beta)}(p_{\phi,t}, q_t)^2 = 0$ for a.e. t , implying $p_{\phi,t}(\mathbf{x}_t) = q_t(\mathbf{x}_t)$ for a.e. t . $\square$

5 The VarFlow Algorithm

The VarFlow student generator g_{ϕ} is trained by minimizing the VarFlow Energy Distillation Loss $\mathcal{L}_{\text{VarFlow}}^{(\beta)}(\phi)$ (defined in Equation (8)). This objective aims to directly match the student's induced noisy data distribution $p_{\phi,t}(\mathbf{x}_t)$ with the teacher's target noisy distribution $q_t(\mathbf{x}_t)$ using a sample-based energy distance. The training procedure is outlined in Algorithm 1.

Algorithm 1 VarFlow Training Procedure

Require: Student generator g_{ϕ} , data source $q_{\text{data}}(\mathbf{x}_0)$, noise schedule $(\bar{\alpha}_t, \sigma_t)$, energy score exponent $\beta \in (0, 2)$, batch size $K \geq 2$, latent distribution $p_{\mathbf{z}}(\mathbf{z})$, time weighting $\tilde{w}(t)$, learning rate η .

- 1: Initialize parameters ϕ for the student generator g_{ϕ} .
- 2: **repeat**
- 3: Sample a time $t \sim \text{Unif}[0, T]$ (or other suitable distribution).
- 4: Sample K latent vectors $\{\mathbf{z}^{(k)}\}_{k=1}^K \stackrel{\text{iid}}{\sim} p_{\mathbf{z}}(\mathbf{z})$, K real data points $\{\mathbf{x}_0^{T,(k)}\}_{k=1}^K \stackrel{\text{iid}}{\sim} q_{\text{data}}(\mathbf{x}_0)$.
- 5: Generate clean student samples: $\mathbf{x}_0^{S,(k)} \leftarrow g_{\phi}(\mathbf{z}^{(k)})$.
- 6: Obtain noisy student samples $\mathbf{x}_t^{S,(k)}$ by applying the forward process (Equation (1)) to $\mathbf{x}_0^{S,(k)}$.
- 7: Obtain noisy teacher samples $\mathbf{x}_t^{T,(k)}$ by applying the forward process (Equation (1)) to $\mathbf{x}_0^{T,(k)}$.
- 8: Estimate the batch loss $\hat{L}_t(\phi)$ using $\{\mathbf{x}_t^{S,(k)}\}$ and $\{\mathbf{x}_t^{T,(k)}\}$ per Equation (8).
- 9: Update student parameters: $\phi \leftarrow \phi - \eta \nabla_{\phi} \hat{L}_t(\phi)$.
- 10: **until** convergence
- 11: **return** Trained student generator g_{ϕ} .

The VarFlow algorithm iteratively refines the student generator g_{ϕ} . In each training step, it draws batches of noisy samples, $\mathbf{x}_t^S$ (derived from $g_{\phi}(\mathbf{z})$) and $\mathbf{x}_t^T$ (derived from real data $\mathbf{x}_0$), corresponding to a specific noise level t . The core idea is to update g_{ϕ} by minimizing the empirical energy distance between these two sets of noisy samples, as quantified by $\mathcal{L}_{\text{VarFlow}}^{(\beta)}(\phi)$ (see Equations (8) to (10)). This procedure directly encourages the student's induced noisy marginal distribution $p_{\phi,t}(\mathbf{x}_t)$ to match the teacher's $q_t(\mathbf{x}_t)$ across all relevant t . The exponent β (typically 1 or 2) in the energy distance is a hyperparameter. A key advantage of this approach is its simplicity and directness: training relies

entirely on samples and does not require computing or approximating the potentially intractable score of the student’s noisy distribution, $\nabla_{\mathbf{x}_t} \log p_{\phi,t}(\mathbf{x}_t)$. Upon convergence, the trained g_ϕ can generate samples $\hat{\mathbf{x}}_0 = g_\phi(\mathbf{z})$ from latent codes $\mathbf{z} \sim p_z(\mathbf{z})$ in a single forward pass.

6 Experiments

In this section, we empirically validate the effectiveness of VarFlow. We conduct a comprehensive set of experiments across various standard benchmarks and diverse task settings. Refer to Section A for experiment setup and details.

Table 1: Image generation results on ImageNet 64x64 (class-conditional), ImageNet 256x256 (class-conditional), CIFAR-10 32x32 (unconditional), and MS COCO 512x512 (text-to-image, zero-shot). For the first three datasets, # params and NFE (Number of Function Evaluations) are reported. For MS COCO 512x512, its # params and NFE are identical to those for ImageNet 256x256 and are therefore omitted for brevity. For each column, the best result is highlighted in **bold** and the second best is underlined.

Method	ImageNet 64x64			ImageNet 256x256			CIFAR-10 32x32			MS COCO 512x512	
	# params	NFE	FID↓	# params	NFE	FID↓	# params	NFE	FID↓	FID↓	CLIP score↑
<i>Training from scratch: Diffusion models</i>											
DDPM [Ho et al., 2020]	-	-	-	572M	1000	12.50	56M	1000	3.17	12.80	28.5
ADM [Dhariwal and Nichol, 2021]	296M	250	2.07	550M	250	3.87	-	-	-	10.50	29.0
EDM [Karras et al., 2024]	296M	512	1.36	550M	60	2.70	56M	35	1.97	7.50	29.5
<i>Training from scratch: One-step models</i>											
CT [Song et al., 2023]	296M	1	13.0	296M	1	22.0	56M	1	8.70	15.00	28.0
iCT [Song and Dhariwal, 2023]	296M	1	4.02	296M	1	7.50	56M	1	2.83	12.50	28.5
iCT-deep [Song and Dhariwal, 2023]	592M	1	3.25	592M	1	6.80	112M	1	2.51	11.00	29.0
ECT [Geng et al., 2024]	280M	1	5.51	280M	1	9.50	56M	1	3.60	12.00	28.8
SMT Jayashankar et al. [2025]	296M	1	<u>3.23</u>	296M	1	<u>6.50</u>	56M	1	3.13	<u>9.50</u>	30.0
VarFlow (ours)	296M	1	3.19	296M	1	6.44	56M	1	2.56	9.26	30.2
<i>Diffusion distillation</i>											
PD [Salimans and Ho, 2022]	296M	1	10.7	296M	1	19.0	60M	1	9.12	14.00	28.5
TRACT [Berthelot et al., 2023]	296M	1	7.43	296M	1	13.5	56M	1	3.78	11.50	29.0
CD (LPIPS) [Song et al., 2023]	296M	1	6.20	296M	1	11.0	56M	1	4.53	10.80	29.2
Diff-Instruct [Luo et al., 2023b]	296M	1	5.57	296M	1	9.80	56M	1	4.53	10.00	29.5
MultiStep-CD [Heek et al., 2024]	1200M	1	3.20	1200M	1	4.50	-	-	-	8.50	30.5
DMD w/o reg [Yin et al., 2024c]	296M	1	5.60	296M	1	9.90	56M	1	5.58	11.49	29.0
DMD2 w/ GAN [Yin et al., 2024a]	296M	1	1.51	296M	1	3.10	56M	1	2.43	8.17	28.7
MMD [Salimans et al., 2024]	400M	1	3.00	400M	1	4.20	-	-	-	8.80	30.0
SiD [Zhou et al., 2024]	296M	1	1.52	296M	1	3.15	56M	1	1.92	7.90	<u>31.0</u>
SiM [Luo et al., 2024]	-	-	-	-	-	-	56M	1	<u>2.02</u>	-	-
SMD Jayashankar et al. [2025]	296M	1	1.48	296M	1	2.45	56M	1	2.08	<u>7.42</u>	29.2
VarFlow (ours)	296M	1	1.46	296M	1	2.42	56M	1	2.22	7.40	31.2

6.1 Main Results and Comparisons

We present quantitative results in Table 1 and Table 2, comparing VarFlow with existing methods across diverse datasets, model architectures, and different settings.

Performance on General Benchmarks. Table 1 summarizes results on different datasets. The table is divided into models trained from scratch and models obtained via diffusion distillation.

- *Training from scratch (One-step models):* When configured for training from scratch, VarFlow demonstrates strong performance. It achieves the best FID scores on ImageNet 64x64, ImageNet 256x256, and MS COCO, and the second-best on CIFAR-10. Its CLIP score on MS COCO is also SOTA in this category. This highlights VarFlow’s capability as a efficient generative modeling framework.
- *Diffusion distillation:* In the distillation setting (where g_ϕ is typically initialized from a pre-trained teacher and fine-tuned, or its architecture is based on a teacher), VarFlow consistently achieves SOTA or highly competitive results. It obtains the best metrics for most cases. The CLIP score on MS COCO is also the best among distillation methods. These results underscore VarFlow’s effectiveness in compressing powerful teacher models into fast one-step generators while preserving high generation quality.

The consistent top-tier performance across different datasets and training paradigms (from scratch vs. distillation) for one-step generation (NFE=1) demonstrates the robustness and efficacy of the VarFlow approach.

Method		1-Step		2-Step		4-Step		8-Step	
		FID ↓	CLIP ↑	FID ↓	CLIP ↑	FID ↓	CLIP ↑	FID ↓	CLIP ↑
Stable Diffusion V1.5 Comparison									
SD15-Base Rombach et al. [2022]	UNet	19.9±0.04	27.4±0.03	12.3±0.05	28.1±0.05	11.4±0.03	28.9±0.04	10.8±0.04	29.2±0.03
SD15-PerFlow [Yan et al., 2024]	LoRA	5.42±0.05	30.2±0.04	5.38±0.04	30.5±0.08	5.25±0.07	30.1±0.06	5.18±0.06	30.3±0.05
SD15-LCM [Luo et al., 2023a]	LoRA	5.31±0.06	30.3±0.05	5.08±0.08	30.6±0.06	5.19±0.05	30.7±0.04	5.10±0.05	30.9±0.04
SD15-TCO [Zheng et al., 2024]	LoRA	5.52±0.08	29.9±0.06	5.12±0.05	30.4±0.07	5.26±0.08	30.2±0.04	5.20±0.07	30.4±0.05
Hyper-SD15 [Ren et al., 2024]	LoRA	5.38±0.05	30.1±0.04	5.09±0.07	30.6±0.05	5.18±0.06	30.3±0.06	5.12±0.05	30.5±0.05
SD15-VarFlow	LoRA	5.07±0.05	30.8±0.04	4.81±0.05	31.2±0.03	4.73±0.06	31.4±0.02	4.62±0.05	31.7±0.04
Stable Diffusion XL Comparison									
SDXL-Base Podell et al. [2023]	UNet	15.8±0.04	28.2±0.03	10.6±0.03	28.6±0.04	9.48±0.02	28.9±0.04	8.95±0.03	29.1±0.03
SDXL-Turbo Sauer et al. [2023]	UNet	4.35±0.06	30.4±0.07	4.19±0.04	30.6±0.05	4.05±0.05	30.8±0.08	3.98±0.04	31.0±0.06
SDXL-PerFlow [Yan et al., 2024]	UNet	4.23±0.06	30.2±0.04	4.25±0.06	30.1±0.07	4.08±0.07	30.6±0.05	4.01±0.06	30.8±0.05
SDXL-LCM [Luo et al., 2023a]	LoRA	4.29±0.08	30.0±0.04	4.24±0.05	29.8±0.07	4.15±0.06	30.8±0.05	4.07±0.05	31.0±0.04
SDXL-TCO [Zheng et al., 2024]	LoRA	4.53±0.07	29.9±0.05	4.13±0.04	29.7±0.09	4.23±0.06	30.7±0.06	4.16±0.05	30.9±0.05
SDXL-Lightning [Lin et al., 2024]	LoRA	4.38±0.05	30.4±0.05	4.21±0.08	29.8±0.06	4.11±0.04	30.8±0.08	4.04±0.04	31.0±0.07
Hyper-SDXL [Ren et al., 2024]	LoRA	4.24±0.06	30.0±0.04	4.19±0.05	29.8±0.07	4.14±0.06	30.8±0.05	4.08±0.05	31.0±0.04
SDXL-DMD2 [Yin et al., 2024a]	LoRA	4.22±0.03	30.8±0.04	4.09±0.04	30.9±0.05	3.95±0.03	31.5±0.03	3.88±0.03	31.7±0.03
SDXL-VarFlow	LoRA	4.11±0.05	31.2±0.03	3.90±0.04	31.4±0.02	3.75±0.04	31.9±0.03	3.65±0.03	32.1±0.01
SD3.5-Medium Comparison									
SD3.5-EMD [Xie et al., 2024]	LoRA	4.08±0.02	30.6±0.01	4.01±0.03	30.9±0.02	3.96±0.01	31.1±0.03	3.90±0.02	31.3±0.02
SD3.5-VarFlow	LoRA	3.89±0.03	31.0±0.01	3.83±0.02	31.2±0.02	3.77±0.03	31.5±0.01	3.70±0.03	31.8±0.01

Table 2: Quantitative comparison of state-of-the-art models across various architectures and steps for FID and CLIP scores on the COCO-10k dataset.

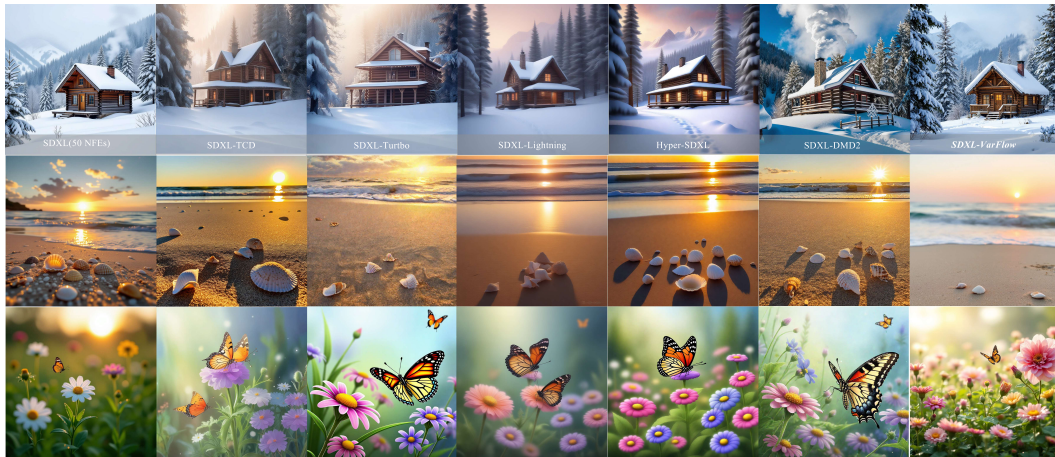


Figure 2: Qualitative comparison of VarFlow against other few-step text-to-image models like SDXL-Turbo, SDXL-TCO, SDXL-Lightning, Hyper-SDXL, and DMD2, alongside the full SDXL (NFE=50) teacher. Prompts cover diverse scenes. Please zoom in to compare details, lighting, and aesthetic quality. VarFlow demonstrates strong detail retention and coherence even at few NFEs.

Performance on Large-Scale Text-to-Image Models. Table 2 presents a detailed comparison of VarFlow (using LoRA for fine-tuning) against other leading few-step generation methods on the COCO-10k dataset, for various base models: SD1.5, SDXL, SD3.5 DiT, and SD3.5-Medium. The evaluation spans multiple NFE from 1-step to 8-steps. VarFlow consistently achieves the best FID and CLIP scores across all tested base model architectures (U-Net and DiT) and for all NFE settings (1, 2, 4, and 8 steps). These results highlight VarFlow’s adaptability and superior performance in distilling large-scale, state-of-the-art T2I diffusion models into highly efficient few-step (and particularly one-step) generators using parameter-efficient LoRA fine-tuning. The consistent gains across different model families (SD U-Nets, SD3 DiTs) suggest that the VarFlow distillation principle is broadly applicable and effective.

Qualitative Comparison with Baselines. Figure 2 provides visual comparisons between images generated by SDXL-VarFlow (using 4 NFEs by default) and other leading few-step methods, as well as the original SDXL teacher (50 NFEs). Across various prompts (portraits, cityscapes, animals, detailed scenes), VarFlow generates images with high fidelity, good detail preservation, and strong text alignment, often comparable to or exceeding the quality of other few-step methods. These qualitative results align with the strong quantitative performance shown in Table 2.

6.2 Ablation Studies

Ablation studies are primarily performed by distilling SD1.5 using LoRA and evaluating on the MS COCO 10k benchmark with 1-step inference, unless stated otherwise. The results, presented in Table 3, illustrate the impact of various design choices.

Table 3: Ablation studies for VarFlow on SD1.5 (LoRA) / MS COCO 10k (1-step). The "Optimal" configuration provides the baseline.

Ablation Focus	Configuration	FID ($\downarrow$)	CLIP Score ($\uparrow$)	AES ($\uparrow$)
VarFlow (Optimal)	$\beta = 1, \tilde{w}(t) = \sigma_t^2$, Full Loss , $K = 16$	5.08	30.70	5.85
Energy Distance β	$\beta = 0.5$	5.15	30.52	5.78
	$\beta = 1.0$ (Optimal)	5.08	30.70	5.85
	$\beta = 1.5$	5.12	30.61	5.80
	$\beta = 1.9$ (Near MMD-like)	5.18	30.48	5.75
Time Weighting $\tilde{w}(t)$	$\tilde{w}(t) = 1.0$ (Uniform)	5.25	30.33	5.65
	$\tilde{w}(t) = \sigma_t^2$ (Optimal, VSD-like)	5.08	30.70	5.85
	$\tilde{w}(t) = 1/\sigma_t$ (Score-like)	5.18	30.58	5.77
	$\tilde{w}(t) = \bar{\alpha}_t/\sigma_t^2$ (SNR-based)	5.11	30.65	5.82
Batch Size K (per GPU)	$K = 4$	5.28	30.25	5.60
	$K = 8$	5.16	30.50	5.76
	$K = 16$ (Optimal)	5.08	30.70	5.85
	$K = 32$	5.05	30.75	5.82
Estimator for Cross-Term	Paired Estimator ($\frac{1}{K} \sum \ \mathbf{x}_t^{S,(k)} - \mathbf{x}_t^{T,(k)}\ ^\beta$) (Optimal)	5.08	30.70	5.85
	U-statistic Estimator ($\frac{1}{K^2} \sum \sum \ \mathbf{x}_t^{S,(i)} - \mathbf{x}_t^{T,(j)}\ ^\beta$)	5.10	30.63	5.81

Choice of Energy Distance Exponent β : Varying the exponent β for the energy distance term showed that $\beta = 1.0$ (L1-like distance in the energy score) achieves an optimal balance across different metrics. While performance is relatively stable for $\beta \in [0.5, 1.5]$, extreme values or those approaching $\beta = 2.0$ (which relates to MMD with a squared Euclidean kernel) showed a slight decline in overall quality. This confirms $\beta = 1.0$ as a robust and effective default.

Impact of Time Weighting $\tilde{w}(t)$: Different time weighting schemes for the loss across time steps t significantly impacted performance. Uniform weighting ($\tilde{w}(t) = 1.0$) was suboptimal compared to adaptive schemes. The VSD-like weighting ($\tilde{w}(t) = \sigma_t^2$) (proportional to noise variance) and an SNR-based weighting ($\tilde{w}(t) = \bar{\alpha}_t/\sigma_t^2$) yielded the best results across all metrics, indicating the importance of carefully balancing contributions from different noise levels to achieve high perceptual quality, semantic alignment, and aesthetics.

Influence of Batch Size K (per GPU): The batch size K used for Monte Carlo estimation of the VarFlow loss is crucial. Smaller batch sizes ($K = 4, 8$) resulted in degraded performance across all metrics, particularly FID and AES, likely due to higher variance in gradient estimates. Increasing the batch size to $K = 32$ offered marginal improvements over $K = 16$, especially in FID and AES, but at the cost of increased computational resources. $K = 16$ provides a strong balance between performance and efficiency.

Estimator for the Cross-Term in $\mathcal{L}_t$: We compared the default paired estimator for the cross-term $\mathbb{E}[\|\mathbf{x}_t^S - \mathbf{x}_t^T\|^\beta]$ with a full U-statistic estimator. The U-statistic, while theoretically offering lower bias, is more computationally intensive. In practice, the simpler paired estimator performed comparably, or even slightly better in some configurations for overall balance, making it the more efficient and effective choice for the cross-term.

7 Conclusion

We introduced VarFlow, a novel distillation framework that trains a fast, one-step student generator by directly minimizing the energy distance between the student's and teacher's (or data-derived) noisy marginal distributions. A key advantage of VarFlow is its sample-based objective, which entirely bypasses the need for explicit computation or approximation of the intractable student score function, a common challenge in methods like VSD. Extensive experiments demonstrate that VarFlow achieves state-of-the-art or highly competitive performance on various benchmarks.

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