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# Image Restoration Through Generalized Ornstein-Uhlenbeck Bridge

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Conghan Yue<sup>1</sup> Zhengwei Peng<sup>1</sup> Junlong Ma<sup>1</sup> Shiyao Du<sup>1</sup> Pengxu Wei<sup>1,2</sup> Dongyu Zhang<sup>1</sup>

## Abstract

Diffusion models exhibit powerful generative capabilities enabling noise mapping to data via reverse stochastic differential equations. However, in image restoration, the focus is on the mapping relationship from low-quality to high-quality images. Regarding this issue, we introduce the Generalized Ornstein-Uhlenbeck Bridge (GOUB) model. By leveraging the natural mean-reverting property of the generalized OU process and further eliminating the variance of its steady-state distribution through the Doob’s  $h$ -transform, we achieve diffusion mappings from point to point enabling the recovery of high-quality images from low-quality ones. Moreover, we unravel the fundamental mathematical essence shared by various bridge models, all of which are special instances of GOUB and empirically demonstrate the optimality of our proposed models. Additionally, we present the corresponding Mean-ODE model adept at capturing both pixel-level details and structural perceptions. Experimental outcomes showcase the state-of-the-art performance achieved by both models across diverse tasks, including inpainting, deraining, and super-resolution. Code is available at <https://github.com/Hammour-steak/GOUB>.

## 1. Introduction

Image restoration involves the restoring of high-quality (HQ) images from their low-quality (LQ) version (Banham & Katsaggelos, 1997; Zhou et al., 1988; Liang et al., 2021; Luo et al., 2023b), which is often characterized as an ill-posed inverse problem due to the loss of crucial information during the degradation from high-quality images to low-quality images. It encompasses a suite of classical tasks, including image deraining (Zhang & Patel, 2017; Yang et al.,

2020; Xiao et al., 2022), denoising (Zhang et al., 2018a; Li et al., 2022; Soh & Cho, 2022; Zhang et al., 2023a), deblurring (Yuan et al., 2007; Kong et al., 2023), inpainting (Jain et al., 2023; Zhang et al., 2023b), and super-resolution (Dong et al., 2015; Zamfir et al., 2023; Wei et al., 2023), among others.

Diffusion models (Sohl-Dickstein et al., 2015; Ho et al., 2020; Song & Ermon, 2019; Song et al., 2021b; Karras et al., 2022) have also been applied to image restoration, yielding favorable results (Ho & Salimans, 2021; Wang et al., 2023; Su et al., 2022; Shi et al., 2024). They mainly follow the standard forward process, diffusing images to pure noise and using low-quality images as conditions to facilitate the generation process of high-quality images (Dhariwal & Nichol, 2021; Ho & Salimans, 2021; Kawar et al., 2021; Saharia et al., 2022; Kawar et al., 2022; Chung et al., 2022b;a; Wang et al., 2023). However, these approaches require the integration of substantial prior knowledge specific to each task such as degradation matrices, limiting their universality.

Furthermore, some studies have attempted to establish a point-to-point mapping from low-quality to high-quality images, learning the general degradation and restoration process and thus circumventing the need for additional prior information for modeling specific tasks (Chen et al., 2022; Cui et al., 2023; Lee et al., 2024). In terms of diffusion models, this mapping can be realized through the bridge (Liu et al., 2022; Su et al., 2022; Liu et al., 2023a), a stochastic process with fixed starting and ending points. By assigning high-quality and low-quality images to the starting and ending points, and initiating with the low-quality images, high-quality images can be obtained by applying the reverse diffusion process, thereby enabling image restoration. However, some bridge models face challenges in learning likelihoods (Liu et al., 2022), necessitating reliance on cumbersome iterative approximation methods (De Bortoli et al., 2021; Su et al., 2022; Shi et al., 2024), which pose significant constraints in practical applications; others do not consider the selection of diffusion process and ignore the optimality of diffusion process (Liu et al., 2023a; Li et al., 2023; Zhou et al., 2024), thus may introducing unnecessary costs and limiting the performance of the model.

This paper proposed a novel image restoration bridge model, the Generalized Ornstein-Uhlenbeck Bridge (GOUB), de-

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<sup>1</sup>Department of Computer Science, Sun Yat-sen University, Guangzhou, Guangdong, China. <sup>2</sup>Pengcheng Laboratory. Correspondence to: Dongyu Zhang <zhangdy27@mail.sysu.edu.cn>.

pictured in Figure 1. Owing to the mean-reverting properties of the Generalized Ornstein-Uhlenbeck (GOU) process, it gradually diffuses the HQ image into a noisy LQ state (denoted as  $\mathbf{x}_T + \lambda\epsilon$  in Figure 1). By applying Doob’s  $h$ -transform on GOU, we modify the diffusion process to eliminate noise on  $\mathbf{x}_T$  to directly bridge the HQ image and its LQ counterpart. The model initiates a point-to-point forward diffusion process and learns its reverse through maximum likelihood estimation, thereby ensuring it can restore a low-quality image to the corresponding high-quality image avoiding the limitation of generality and costly iterative approximation. Our main contributions can be summarized as follows:

- We introduce a novel image restoration bridge model GOUB which eliminates variance of the ending point on the GOU process, directly connecting the high and low-quality images and is particularly expressive in deep visual features and diversity.
- Benefiting from the distinctive features of the parameterization mechanism, we introduce the corresponding Mean-ODE model, demonstrating a strong ability to capture pixel-level details and structural perceptions.
- We uncover the mathematical essence of several bridge models, all of which are special cases of the GOUB, and empirically demonstrate the optimality of our proposed models.
- Our model has achieved state-of-the-art results on numerous image restoration tasks, such as inpainting, deraining, and super-resolution.

## 2. Preliminaries

### 2.1. Score-based Diffusion Model

The score-based diffusion model (Sohl-Dickstein et al., 2015; Ho et al., 2020; Song et al., 2021b) is a category of generative model that seamlessly transitions data into noise via a diffusion process and generates samples by learning and adapting the reverse process (Anderson, 1982). Assuming a dataset consists of  $n$  dimensional independent identically distributed (i.i.d.) samples, following an unknown distribution denoted by  $p(\mathbf{x}_0)$ . The time-dependent forward process of the diffusion model can be described by the following SDE:

$$d\mathbf{x}_t = \mathbf{f}(\mathbf{x}_t, t) dt + g_t d\mathbf{w}_t, \quad (1)$$

where  $\mathbf{f} : \mathbb{R}^n \rightarrow \mathbb{R}^n$  is the drift coefficient,  $g_t : \mathbb{R} \rightarrow \mathbb{R}$  is the scalar diffusion coefficient and  $\mathbf{w}_t$  denotes the standard Brownian motion. Typically,  $p(\mathbf{x}_0)$  evolves over time  $t$  from 0 to a sufficiently large  $T$  into  $p(\mathbf{x}_T)$  through the

SDE, such that  $p(\mathbf{x}_T)$  will approximate a standard Gaussian distribution  $p_{\text{prior}}(\mathbf{x})$ . Meanwhile, the forward SDE has a corresponding reverse time SDE (Anderson, 1982) whose closed form is given by:

$$d\mathbf{x}_t = [\mathbf{f}(\mathbf{x}_t, t) - g_t^2 \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t)] dt + g_t d\mathbf{w}_t. \quad (2)$$

Starting from time  $T$ ,  $p(\mathbf{x}_T)$  can progressively transform to  $p(\mathbf{x}_0)$  by traversing the trajectory of the reverse SDE. The score  $\nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t)$  can generally be parameterized as  $\mathbf{s}_\theta(\mathbf{x}_t, t)$  and employ conditional score matching (Vincent, 2011) as the loss function for training:

$$\begin{aligned} \mathcal{L} &= \frac{1}{2} \int_0^T \mathbb{E}_{\mathbf{x}_t} \left[ \lambda(t) \|\nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t) - \mathbf{s}_\theta(\mathbf{x}_t, t)\|^2 \right] dt \\ &\propto \frac{1}{2} \int_0^T \mathbb{E}_{\mathbf{x}_0, \mathbf{x}_t} \left[ \lambda(t) \|\nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t | \mathbf{x}_0) - \mathbf{s}_\theta(\mathbf{x}_t, t)\|^2 \right] dt, \end{aligned} \quad (3)$$

where  $\lambda(t)$  serves as a weighting function, and if selected as  $g_t^2$  that yields a more optimal upper bound on the negative log-likelihood (Song et al., 2021a). The second line is actually the most commonly used, as the conditional probability  $p(\mathbf{x}_t | \mathbf{x}_0)$  is generally accessible. Ultimately, one can sample  $\mathbf{x}_T$  from the prior distribution  $p(\mathbf{x}_T) \approx p_{\text{prior}}(\mathbf{x})$  and obtain the  $\mathbf{x}_0$  through the numerical solution of Equation (2) via iterative steps, thereby completing the generation process.

### 2.2. Generalized Ornstein-Uhlenbeck process

The Generalized Ornstein-Uhlenbeck (GOU) process is the time-varying OU process (Ahmad, 1988). It is a stationary Gaussian-Markov process, whose marginal distribution gradually tends towards a stable mean and variance over time. The GOU process is generally defined as follows:

$$d\mathbf{x}_t = \theta_t (\boldsymbol{\mu} - \mathbf{x}_t) dt + g_t d\mathbf{w}_t, \quad (4)$$

where  $\boldsymbol{\mu}$  is a given state vector,  $\theta_t$  denotes a scalar drift coefficient and  $g_t$  represents the diffusion coefficient. At the same time, we require  $\theta_t, g_t$  to satisfy the specified relationship  $2\lambda^2 = g_t^2/\theta_t$ , where  $\lambda^2$  is a given constant scalar. As a result, its transition probability possesses a closed-form analytical solution:

$$\begin{aligned} p(\mathbf{x}_t | \mathbf{x}_s) &= N(\bar{\mathbf{m}}_{s:t}, \bar{\sigma}_{s:t}^2 \mathbf{I}) = \\ &N\left(\boldsymbol{\mu} + (\mathbf{x}_s - \boldsymbol{\mu}) e^{-\bar{\theta}_{s:t}}, \frac{g_t^2}{2\theta_t} \left(1 - e^{-2\bar{\theta}_{s:t}}\right) \mathbf{I}\right), \quad (5) \\ \bar{\theta}_{s:t} &= \int_s^t \theta_z dz. \end{aligned}$$

A simple proof is provided in Appendix C. For the sake of simplicity in subsequent representations, we denote  $\bar{\theta}_{0:t}$  and  $\bar{\sigma}_{0:t}$  as  $\bar{\theta}_t$  and  $\bar{\sigma}_t$  respectively. Consequently,  $p(\mathbf{x}_t)$  will steadily converge towards a Gaussian distribution with the mean of  $\boldsymbol{\mu}$  and the variance of  $\lambda^2$  as time  $t$  progresses meaning that it exhibits the mean-reverting property.

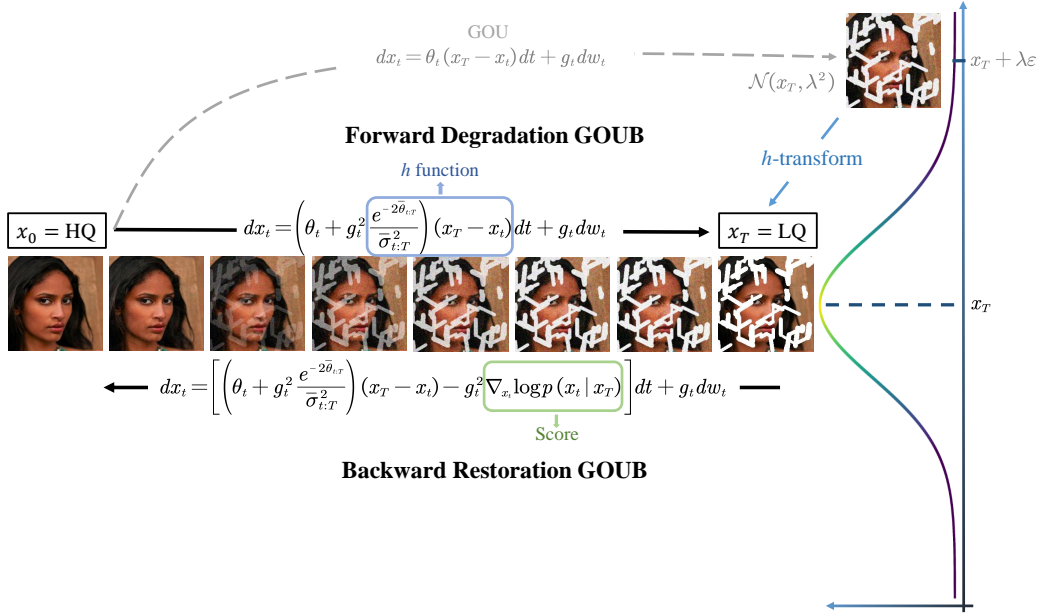


Figure 1. Overview of the proposed GOUB for image restoration. The GOU process is capable of transferring an HQ image into a noisy LQ image. Additionally, through the application of  $h$ -transform, we can eliminate the noise on LQ, enabling the GOUB model to precisely bridge the gap between HQ and LQ.

### 2.3. Doob’s $h$ -transform

Doob’s  $h$ -transform (Särkkä & Solin, 2019) is a mathematical technique applied to stochastic processes. It involves transforming the original process by incorporating a specific  $h$ -function into the drift term of the SDE, modifying the process to pass through a predetermined terminal point. More precisely, given the SDE (1), if it is desired to pass through the given fixed point  $\mathbf{x}_T$  at  $t = T$ , an additional drift term must be incorporated into the original SDE:

$$d\mathbf{x}_t = [\mathbf{f}(\mathbf{x}_t, t) + g_t^2 \mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T)] dt + g_t d\mathbf{w}_t, \quad (6)$$

where  $\mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T) = \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_T | \mathbf{x}_t)$  and  $\mathbf{x}_0$  starts from  $p(\mathbf{x}_0 | \mathbf{x}_T)$ . A simple proof can be found in Appendix D. In comparison to (1), the marginal distribution of (6) is conditioned on  $\mathbf{x}_T$ , with its forward conditional probability density given by  $p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)$  satisfying the forward Kolmogorov equation that is defined by (6). Intuitively,  $p(\mathbf{x}_T | \mathbf{x}_0, \mathbf{x}_T) = 1$  at  $t = T$ , ensuring that the SDE invariably passes through the specified point  $\mathbf{x}_T$  for any initial state  $\mathbf{x}_0$ .

## 3. GOUB

The GOU process (4) is characterized by mean-reverting properties that if we consider the initial state  $\mathbf{x}_0$  to represent a high-quality image and the corresponding low-quality image  $\mathbf{x}_T = \boldsymbol{\mu}$  as the final condition, then the high-quality image will gradually converge to a Gaussian distribution with the low-quality image as its mean and a stable variance

$\lambda^2$ . This naturally connects some information between high and low-quality images, offering an inherent advantage in image restoration. However, the initial state of the reverse process necessitates the artificial addition of noise to low-quality images, resulting in certain information loss and thus affecting the performance (Luo et al., 2023a).

In actuality, we are more focused on the connections between points (Liu et al., 2022; De Bortoli et al., 2021; Su et al., 2022; Li et al., 2023; Zhou et al., 2024) in image restoration. Coincidentally, the Doob’s  $h$ -transform technique can modify an SDE such that it passes through a specified  $\mathbf{x}_T$  at terminal time  $T$ . Accordingly, it is crucial to note that the application of the  $h$ -transform to the GOU process effectively eliminates the impact of terminal noise, directly bridging a point-to-point relationship between high-quality and low-quality images.

### 3.1. Forward and backward process

Applying the  $h$ -transform, we can readily derive the forward process of the GOUB, leading to the following proposition:

**Proposition 3.1.** *Let  $\mathbf{x}_t$  be a finite random variable describing by the given generalized Ornstein-Uhlenbeck process (4), suppose  $\mathbf{x}_T = \boldsymbol{\mu}$ , the evolution of its marginal distribution  $p(\mathbf{x}_t | \mathbf{x}_T)$  satisfies the following SDE:*

$$d\mathbf{x}_t = \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t,T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) dt + g_t d\mathbf{w}_t. \quad (7)$$

Additionally, the forward transition  $p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)$  is given

by:

$$\begin{aligned}
 p(\mathbf{x}_t \mid \mathbf{x}_0, \mathbf{x}_T) &= N(\bar{\mathbf{m}}'_t, \bar{\sigma}_t'^2 \mathbf{I}), \\
 \bar{\mathbf{m}}'_t &= e^{-\bar{\theta}_t} \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{x}_0 + \left[ \left(1 - e^{-\bar{\theta}_t}\right) \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} + e^{-2\bar{\theta}_{t:T}} \frac{\bar{\sigma}_t^2}{\bar{\sigma}_T^2} \right] \mathbf{x}_T \\
 \bar{\sigma}_t'^2 &= \frac{\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2}
 \end{aligned} \tag{8}$$

The derivation of the proposition is provided in the Appendix A.1. With Proposition 3.1, there is no need to perform multi-step forward iteration using the SDE; instead, we can directly use its closed-form solution for one-step forward sampling.

Similarly, applying the previous SDE theory enables us to easily derive the reverse process, which leads to the following Proposition 3.2:

**Proposition 3.2.** *The reverse SDE of equation (7) has a marginal distribution  $p(\mathbf{x}_t \mid \mathbf{x}_T)$ , and is given by:*

$$\begin{aligned}
 d\mathbf{x}_t &= \left[ \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) \right. \\
 &\quad \left. - g_t^2 \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t \mid \mathbf{x}_T) \right] dt + g_t d\mathbf{w}_t,
 \end{aligned} \tag{9}$$

and exists a probability flow ODE:

$$\begin{aligned}
 d\mathbf{x}_t &= \left[ \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) \right. \\
 &\quad \left. - \frac{1}{2} g_t^2 \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t \mid \mathbf{x}_T) \right] dt.
 \end{aligned} \tag{10}$$

We are capable of initiating from a low-quality image  $\mathbf{x}_T$  and proceeding to utilize Euler sampling solving the reverse SDE or ODE for restoration purposes.

### 3.2. Training object

The score term  $\nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t \mid \mathbf{x}_T)$  can be parameterized by a neural network  $\mathbf{s}_\theta(\mathbf{x}_t, \mathbf{x}_T, t)$  and can be estimated using the loss function (3). Unfortunately, training the score function for SDEs generally presents a significant challenge. Nevertheless, since the analytical form of GOUB is directly obtainable, we will introduce the use of maximum likelihood for training, which yields a more stable loss function.

We first discretize the continuous time interval  $[0, T]$  into  $N$  sufficiently fine-grained intervals in a reasonable manner, denoted as  $\{\mathbf{x}_t\}_{t \in [0, N]}$ ,  $\mathbf{x}_N = \mathbf{x}_T$ . We are concerned with maximizing the log-likelihood, which leads us to the following proposition:

**Proposition 3.3.** *Let  $\mathbf{x}_t$  be a finite random variable describing by the given generalized Ornstein-Uhlenbeck process (4), for a fixed  $\mathbf{x}_T$ , the expectation of log-likelihood  $\mathbb{E}_{p(\mathbf{x}_0)}[\log p_\theta(\mathbf{x}_0 \mid \mathbf{x}_T)]$  possesses an Evidence Lower Bound (ELBO):*

$$\begin{aligned}
 ELBO &= \mathbb{E}_{p(\mathbf{x}_0)} \left[ \mathbb{E}_{p(\mathbf{x}_1 \mid \mathbf{x}_0)} [\log p_\theta(\mathbf{x}_0 \mid \mathbf{x}_1, \mathbf{x}_T)] - \right. \\
 &\quad \left. \sum_{t=2}^T \mathbb{E}_{p(\mathbf{x}_t \mid \mathbf{x}_0)} [KL(p(\mathbf{x}_{t-1} \mid \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T) \parallel p_\theta(\mathbf{x}_{t-1} \mid \mathbf{x}_t, \mathbf{x}_T))] \right]
 \end{aligned} \tag{11}$$

Assuming  $p_\theta(\mathbf{x}_{t-1} \mid \mathbf{x}_t, \mathbf{x}_T)$  is a Gaussian distribution with a constant variance  $N(\boldsymbol{\mu}_{\theta, t-1}, \sigma_{\theta, t-1}^2 \mathbf{I})$ , maximizing the ELBO is equivalent to minimizing:

$$\mathcal{L} = \mathbb{E}_{t, \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T} \left[ \frac{1}{2\sigma_{\theta, t-1}^2} \|\boldsymbol{\mu}_{t-1} - \boldsymbol{\mu}_{\theta, t-1}\|^2 \right], \tag{12}$$

where  $\boldsymbol{\mu}_{t-1}$  represents the mean of  $p(\mathbf{x}_{t-1} \mid \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T)$ :

$$\boldsymbol{\mu}_{t-1} = \frac{1}{\bar{\sigma}_t'^2} [\bar{\sigma}_{t-1}'^2 (\mathbf{x}_t - b\mathbf{x}_T) a + (\bar{\sigma}_t'^2 - \bar{\sigma}_{t-1}'^2 a^2) \bar{\mathbf{m}}'_t], \tag{13}$$

where,

$$\begin{aligned}
 a &= \frac{e^{-\bar{\theta}_{t-1:T}} \bar{\sigma}_{t:T}^2}{\bar{\sigma}_{t-1:T}^2}, \\
 b &= \frac{1}{\bar{\sigma}_T^2} \left\{ (1 - e^{-\bar{\theta}_t}) \bar{\sigma}_{t:T}^2 + e^{-2\bar{\theta}_{t:T}} \bar{\sigma}_t^2 \right. \\
 &\quad \left. - \left[ (1 - e^{-\bar{\theta}_{t-1}}) \bar{\sigma}_{t-1:T}^2 + e^{-2\bar{\theta}_{t-1:T}} \bar{\sigma}_{t-1}^2 \right] a \right\}
 \end{aligned}$$

The derivation of the proposition is provided in the Appendix A.2. With Proposition 3.3, we can easily construct the training objective. In this work, we try to parameterized  $\boldsymbol{\mu}_{\theta, t-1}$  from differential of SDE which can be derived from equation (9):

$$\begin{aligned}
 \mathbf{x}_{t-1} &= \mathbf{x}_t - \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) \\
 &\quad + g_t^2 \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t \mid \mathbf{x}_T) - g_t \boldsymbol{\epsilon}_t,
 \end{aligned} \tag{14}$$

where  $\boldsymbol{\epsilon}_t \sim N(\mathbf{0}, dt\mathbf{I})$ , therefore:

$$\begin{aligned}
 \boldsymbol{\mu}_{\theta, t-1} &= \mathbf{x}_t - \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) \\
 &\quad + g_t^2 \nabla_{\mathbf{x}_t} \log p_\theta(\mathbf{x}_t \mid \mathbf{x}_T), \\
 \sigma_{\theta, t-1} &= g_t.
 \end{aligned} \tag{15}$$

Inspired by conditional score matching, we can parameterize noise as  $\boldsymbol{\epsilon}_\theta(\mathbf{x}_t, \mathbf{x}_T, t)$ , thus the score  $\nabla_{\mathbf{x}_t} \log p_\theta(\mathbf{x}_t \mid \mathbf{x}_T)$  can be represented as  $-\boldsymbol{\epsilon}_\theta(\mathbf{x}_t, \mathbf{x}_T, t) / \bar{\sigma}_t'$ . In addition, during our empirical research, we found that utilizing L1 loss

yields enhanced image reconstruction outcomes (Boyd & Vandenberghe, 2004; Hastie et al., 2009). This approach enables the model to learn pixel-level details more easily, resulting in markedly improved visual quality. Therefore, the final training object is:

$$\begin{aligned} \mathcal{L} = \mathbb{E}_{t, \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T} & \left[ \frac{1}{2g_t^2} \left\| \frac{1}{\bar{\sigma}_t'^2} [\bar{\sigma}_{t-1}'^2 (\mathbf{x}_t - b\mathbf{x}_T) a \right. \right. \\ & \left. \left. + (\bar{\sigma}_t'^2 - \bar{\sigma}_{t-1}'^2 a^2) \bar{\mathbf{m}}_t' \right] - \mathbf{x}_t \right. \\ & \left. + \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) \right. \\ & \left. + \frac{g_t^2}{\bar{\sigma}_t} \epsilon_{\theta}(\mathbf{x}_t, \mathbf{x}_T, t) \right\| \right] \end{aligned} \quad (16)$$

Consequently, if we obtain the optimal  $\epsilon_{\theta}^*(\mathbf{x}_t, \mathbf{x}_T, t)$ , we can compute the score  $\nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t | \mathbf{x}_T) \approx -\epsilon_{\theta}^*(\mathbf{x}_t, \mathbf{x}_T, t) / \bar{\sigma}_t'$  for reverse process. Starting from a low-quality image  $\mathbf{x}_T$ , we can recover  $\mathbf{x}_0$  by using Equation (9) to perform reverse iteration.

### 3.3. Mean-ODE

Unlike normal diffusion models, our parameterization of the mean  $\mu_{\theta, t-1}$  is derived from the differential of SDE which effectively combines the characteristics of discrete diffusion models and continuous score-based generative models. In the reverse process, the value of each sampling step will be approximated to the true mean during training. Therefore, we propose a Mean-ODE model, which omits the Brownian drift term:

$$\begin{aligned} d\mathbf{x}_t = & \left[ \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) \right. \\ & \left. - g_t^2 \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t | \mathbf{x}_T) \right] dt, \end{aligned} \quad (17)$$

To simplify the expression, we use GOUB to represent the GOUB (SDE) sampling model and Mean-ODE to represent the GOUB (Mean-ODE) sampling model. Our following experiments have demonstrated that the Mean-ODE is more effective than the corresponding Score-ODE at capturing the pixel details and structural perceptions of images, playing a pivotal role in image restoration tasks. Concurrently, the SDE model (9) is more focused on deep visual features and diversity.

## 4. Experiments

We conduct experiments under three popular image restoration tasks: image inpainting, image deraining, and image super-resolution. Four metrics are employed for the model

Table 1. **Image Inpainting.** Qualitative comparison with the relevant baselines on CelebA-HQ.

METHOD	PSNR $\uparrow$	SSIM $\uparrow$	LPIPS $\downarrow$	FID $\downarrow$
PromptIR	30.22	0.9180	0.068	32.69
DDRM	27.16	0.8993	0.089	37.02
IR-SDE	28.37	0.9166	0.046	25.13
GOUB	28.98	0.9067	<b>0.037</b>	<b>4.30</b>
Mean-ODE	<b>31.39</b>	<b>0.9392</b>	0.052	12.24

evaluation, *i.e.*, Peak Signal-to-Noise Ratio (PSNR) for assessing reconstruction quality, Structural Similarity Index (SSIM) (Wang et al., 2004) for gauging structural perception, Learned Perceptual Image Patch Similarity (LPIPS) (Zhang et al., 2018b) for evaluating the depth and quality of features, and Fréchet Inception Distance (FID) (Heusel et al., 2017) to measure the diversity in generated images. More experiment details are present in Appendix E.

**Image Inpainting.** Image inpainting involves filling in missing or damaged parts of an image, to restore or enhance the overall visual effect of the image. We have selected the CelebA-HQ  $256 \times 256$  datasets (Karras et al., 2018) for both training and testing with 100 thin masks. We compare our models with several current baseline inpainting approaches such as PromptIR (Potlapalli et al., 2023), DDRM (Kawar et al., 2022) and IR-SDE (Luo et al., 2023a). The relevant experimental results are shown in the Table 1 and Figure 2. It is observed that the two proposed models achieved state-of-the-art results in their respective areas of strength and also delivered highly competitive outcomes on other metrics. From a visual perspective, our model excels in capturing details such as eyebrows, eyes, and image backgrounds.

**Image Deraining.** We have selected the Rain100H datasets (Yang et al., 2017) for our training and testing, which includes 1800 pairs of training data and 100 images for testing. It is important to note that in this task, similar to other deraining models, we present the PSNR and SSIM scores specifically on the Y channel (YCbCr space). We report state-of-the-art approaches for comparison: MPRNet (Zamir et al., 2021), M3SNet-32 (Gao et al., 2023), MAXIM (Tu et al., 2022), MHNet (Gao & Dang, 2023), IR-SDE (Luo et al., 2023a). The relevant experimental results are shown in the Table 2 and Figure 3. Similarly, both models achieved SOTA results respectively in the deraining task. Visually, it can be also observed that our model excels in capturing details such as the moon, the sun, and tree branches.

**Image Super-Resolution.** Single image super-resolution aims to recover a higher resolution and clearer version from a low-resolution image. We conducted training and evalua-

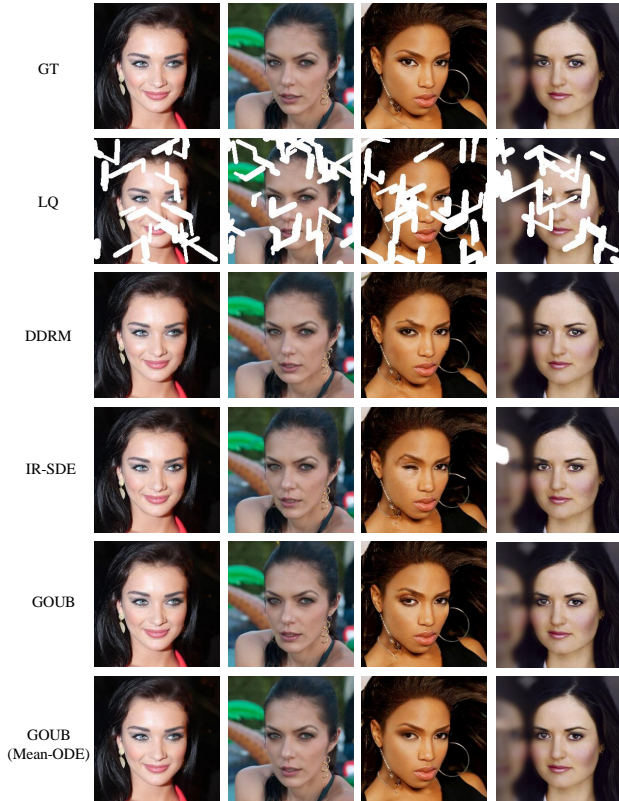


Figure 2. Qualitative comparison of the visual results of different inpainting methods on the CelebA-HQ dataset with thin mask.

tion on the DIV2K validation set for  $4\times$  upscaling (Agustsson & Timofte, 2017) and all low-resolution images were bicubically rescaled to the same size as their corresponding high-resolution images. To show that our models are in line with the state-of-the-art, we compare to the DDRM (Kawar et al., 2022) and IR-SDE (Luo et al., 2023a). The relevant experimental results are provided in Table 3 and Figure 4. As can be seen, our GOUB is superior to benchmarks in various indicators and handles visual details better such as edges and hair.

**Superiority of Mean-ODE.** Additionally, we conduct ablation experiments using the corresponding Score-ODE (10) model to demonstrate the superiority of our proposed Mean-ODE model in image restoration. From Table 4, it is evident that the performance of Mean-ODE is significantly superior to that of the corresponding Score-ODE. This is because the sampling results of each sampling step of Mean-ODE directly approximate the true mean during the training process, as opposed to the parameterized approach such as DDPM, which relies on expectations. Consequently, our proposed Mean-ODE demonstrates better reconstruction effects and is more suitable for image restoration tasks.

Table 2. **Image Deraining.** Qualitative comparison with the relevant baselines on Rain100H.

METHOD	PSNR $\uparrow$	SSIM $\uparrow$	LPIPS $\downarrow$	FID $\downarrow$
MPRNet	30.41	0.8906	0.158	61.59
M3SNet-32	30.64	0.8920	0.154	60.26
MAXIM	30.81	0.9027	0.133	58.72
MHNet	31.08	0.8990	0.126	57.93
IR-SDE	31.65	0.9041	0.047	18.64
GOUB	31.96	0.9028	<b>0.046</b>	<b>18.14</b>
Mean-ODE	<b>34.56</b>	<b>0.9414</b>	0.077	32.83

Table 3. **Image  $4\times$  Super-Resolution.** Qualitative comparison with the relevant baselines on DIV2K.

METHOD	PSNR $\uparrow$	SSIM $\uparrow$	LPIPS $\downarrow$	FID $\downarrow$
DDRM	24.35	0.5927	0.364	78.71
IR-SDE	25.90	0.6570	0.231	45.36
GOUB	26.89	0.7478	<b>0.220</b>	<b>20.85</b>
Mean-ODE	<b>28.50</b>	<b>0.8070</b>	0.328	22.14

## 5. Analysis

The Doob’s  $h$ -transform of the generalized Ornstein-Uhlenbeck process, also known as the conditional GOU process has been an intriguing topic in previous applied mathematical research (Salminen, 1984; Cheridito et al., 2003; Heng et al., 2021). On account of the mean-reverting property of the GOU process, applying the  $h$ -transform makes it most straightforward to eliminate the variance and drive it towards a Dirac distribution in its steady state which is highly advantageous for its applications in image restoration. In previous research on diffusion models, there has been limited focus on the cases of  $f$  or  $g$ , and generally used the VE process (Song et al., 2021b) represented by NCSN (Song & Ermon, 2019) or the VP process (Song et al., 2021b) represented by DDPM (Ho et al., 2020).

In this section, we demonstrate that the mathematical essence of several recent meaningful diffusion bridge models is the same (Li et al., 2023; Zhou et al., 2024; Liu et al., 2023a) and they all represent Brownian bridge (Chow, 2009) models, details are provided in the Appendix B.1. Then, we also found that the VE and VP processes are special cases of GOU, leading to the following proposition:

**Proposition 5.1.** *For a given GOU process (4), there exists relationships:*

$$\begin{aligned} \lim_{\theta_t \rightarrow 0} \text{GOU} &= \text{VE} \\ \lim_{\mu \rightarrow 0, \lambda \rightarrow 1} \text{GOU} &= \text{VP} \end{aligned} \quad (18)$$

Details are provided in the Appendix B.2. Therefore, we

Table 4. Qualitative comparison with the corresponding Score-ODE on various tasks.

METHOD	Image Inpainting				Image Deraining				Image 4× Super-Resolution			
	PSNR↑	SSIM↑	LPIPS↓	FID↓	PSNR↑	SSIM↑	LPIPS↓	FID↓	PSNR↑	SSIM↑	LPIPS↓	FID↓
Score-ODE	18.23	0.6266	0.389	161.54	13.64	0.7404	0.338	191.15	28.14	0.7993	0.344	25.51
Mean-ODE	<b>31.39</b>	<b>0.9392</b>	<b>0.052</b>	<b>12.24</b>	<b>34.56</b>	<b>0.9414</b>	<b>0.077</b>	<b>32.83</b>	<b>28.50</b>	<b>0.8070</b>	<b>0.328</b>	<b>22.14</b>

Table 5. Qualitative comparison with the different bridge models on CelebA-HQ, Rain100H, and DIV2K datasets.

METHOD	Image Inpainting				Image Deraining				Image 4× Super-Resolution			
	PSNR↑	SSIM↑	LPIPS↓	FID↓	PSNR↑	SSIM↑	LPIPS↓	FID↓	PSNR↑	SSIM↑	LPIPS↓	FID↓
VEB	27.75	0.8943	0.056	13.70	30.39	0.8975	0.059	28.54	24.21	0.5808	0.384	36.55
VPB	27.32	0.8841	0.049	11.87	30.89	0.8847	0.051	23.36	25.40	0.6041	0.342	29.17
GOUB	<b>28.98</b>	<b>0.9067</b>	<b>0.037</b>	<b>4.30</b>	<b>31.96</b>	<b>0.9028</b>	<b>0.046</b>	<b>18.14</b>	<b>26.89</b>	<b>0.7478</b>	<b>0.220</b>	<b>20.85</b>

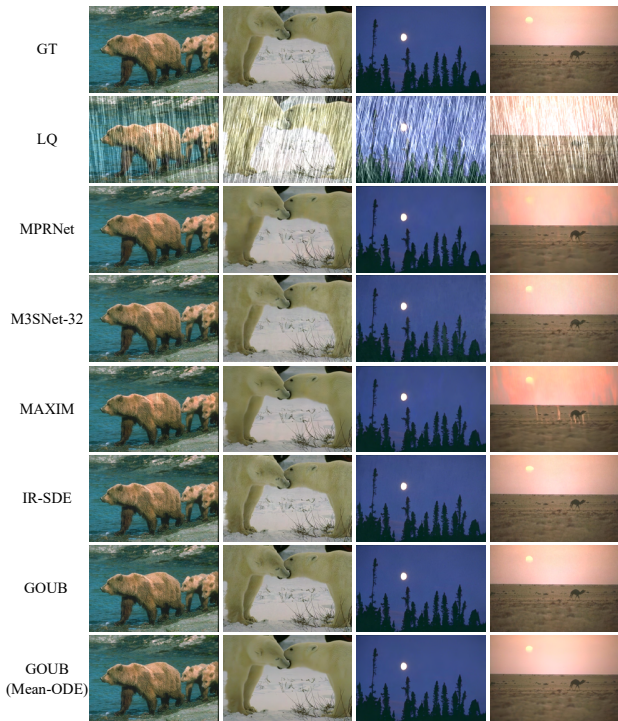


Figure 3. Qualitative comparison of the visual results of different deraining methods on the Rain100H dataset.

conduct experiments on VE Bridge (VEB) (Li et al., 2023; Zhou et al., 2024; Liu et al., 2023a) and VP Bridge (VPB) (Zhou et al., 2024) to demonstrate the optimality of our proposed GOUB model in image restoration. We keep all the model hyperparameters consistent and results are shown in Table 5 and Figure 5.

It can be seen that under the same configuration of model hyperparameters, the performance of the GOUB is notably superior to the other two types of bridge models, which

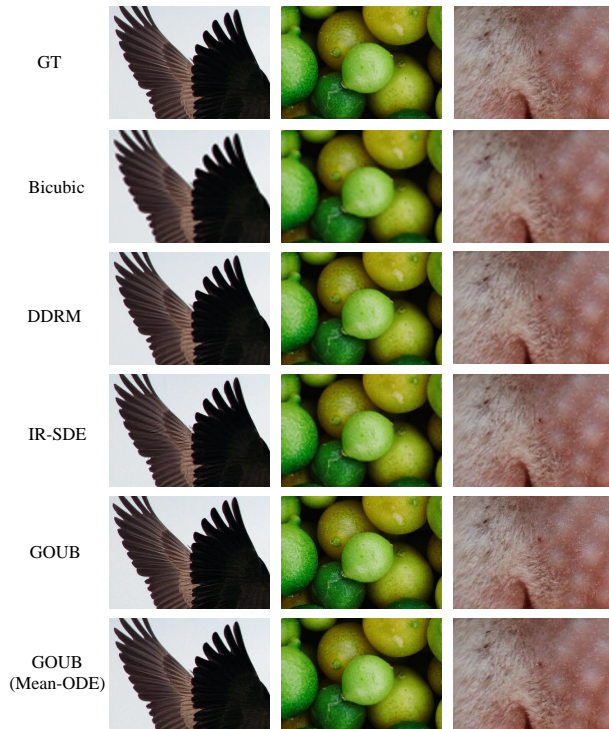


Figure 4. Qualitative comparison of the visual results of different 4x super-resolution methods on the DIV2K dataset.

demonstrates the optimality of GOUB and also highlights the importance of the choice of diffusion process in diffusion models.

## 6. Related Works

**Conditional Generation.** As previously highlighted, in the work of image restoration using diffusion models, the focus of some research has predominantly been on using low-

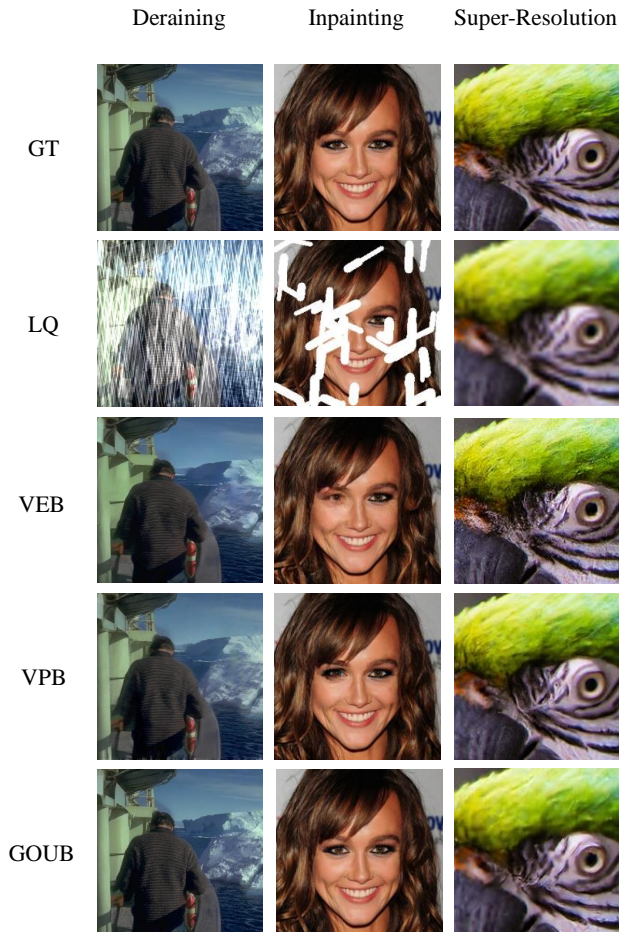


Figure 5. Qualitative comparison with the different bridge models in many tasks.

quality images as conditional inputs  $y$  to guide the generation process. They (Kawar et al., 2021; Saharia et al., 2022; Kawar et al., 2022; Chung et al., 2022a;b; 2023; Zhao et al., 2023; Murata et al., 2023; Feng et al., 2023) all endeavor to solve or approximate the classifier  $\log \nabla_{\mathbf{x}_t} p(\mathbf{y} | \mathbf{x}_t)$ , necessitating the incorporation of additional prior knowledge to model specific degradation processes which both complex and lacking in universality.

**Diffusion Bridge.** This segment of work obviates the need for prior knowledge, constructing a diffusion bridge model from high-quality to low-quality images, thereby learning the degradation process. The previously mentioned approach (Liu et al., 2022; De Bortoli et al., 2021; Su et al., 2022; Liu et al., 2023a; Shi et al., 2024; Li et al., 2023; Zhou et al., 2024; Albergo et al., 2023) fall into this class and are characterized by the issues of significant computational expense in solution seeking and also not the optimal model framework. Additionally, some models of flow category (Lipman et al., 2023; Liu et al., 2023b; Tong et al., 2023;

Albergo & Vanden-Eijnden, 2023; Delbracio & Milanfar, 2023) also belong to the diffusion bridge models and face the similar issue.

## 7. Conclusion

In this paper, we introduced the Generalized Ornstein-Uhlenbeck Bridge (GOUB) model, a diffusion bridge model that applies the Doob’s  $h$ -transform to the GOU process. This model can address general image restoration tasks without the need for specific prior knowledge. Furthermore, we have uncovered the mathematical essence of several bridge models and empirically demonstrated the optimality of our proposed model. In addition, considering our unique mean parameterization mechanism, we proposed the Mean-ODE model. Experimental results indicate that both models achieve state-of-the-art results in their respective areas of strength on various tasks, including inpainting, deraining, and super-resolution. We believe that the exploration of diffusion process and bridge models holds significant importance not only in the field of image restoration but also in advancing the study of generative diffusion models.

## Acknowledgements

This work is supported in part by National Natural Science Foundation of China (NSFC) under Grant No. 62376292, U21A20470, and Guangdong Basic and Applied Basic Research Foundation under Grant No. 2024A1515011741.

## Impact Statement

This paper presents work whose goal is to advance the field of Machine Learning. There are many potential societal consequences of our work, none which we feel must be specifically highlighted here.

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## A. Proof

### A.1. Proof of Proposition 3.1

**Proposition 3.1.** Let  $\mathbf{x}_t$  be a finite random variable describing by the given generalized Ornstein-Uhlenbeck process (4), suppose  $\mathbf{x}_T = \boldsymbol{\mu}$ , the evolution of its marginal distribution  $p(\mathbf{x}_t | \mathbf{x}_T)$  satisfies the following SDE:

$$d\mathbf{x}_t = \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) dt + g_t d\mathbf{w}_t, \quad (7)$$

additionally, the forward transition  $p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)$  is given by:

$$\begin{aligned} p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T) &= N(\bar{\mathbf{m}}'_t, \bar{\sigma}_t'^2 \mathbf{I}) \\ &= N\left( e^{-\bar{\theta}_t} \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{x}_0 + \left[ (1 - e^{-\bar{\theta}_t}) \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} + e^{-2\bar{\theta}_{t:T}} \frac{\bar{\sigma}_t^2}{\bar{\sigma}_T^2} \right] \mathbf{x}_T, \frac{\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{I} \right) \end{aligned} \quad (8)$$

*Proof:* Based on (5), we have:

$$p(\mathbf{x}_t | \mathbf{x}_0) = N\left( \mathbf{x}_T + (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_t}, \bar{\sigma}_t^2 \mathbf{I} \right) \quad (19)$$

$$p(\mathbf{x}_T | \mathbf{x}_t) = N\left( \mathbf{x}_T + (\mathbf{x}_t - \mathbf{x}_T) e^{-\bar{\theta}_{t:T}}, \bar{\sigma}_{t:T}^2 \mathbf{I} \right) \quad (20)$$

$$p(\mathbf{x}_T | \mathbf{x}_0) = N\left( \mathbf{x}_T + (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_T}, \bar{\sigma}_T^2 \mathbf{I} \right) \quad (21)$$

Firstly, the  $h$  function can be directly compute:

$$\begin{aligned} \mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T) &= \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_T | \mathbf{x}_t) \\ &= -\nabla_{\mathbf{x}_t} \frac{(\mathbf{x}_t - \mathbf{x}_T)^2 e^{-2\bar{\theta}_{t:T}}}{2\bar{\sigma}_{t:T}^2} \\ &= (\mathbf{x}_T - \mathbf{x}_t) \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \end{aligned} \quad (22)$$

Therefore, followed by Doob's  $h$ -transform (6), the SDE of marginal distribution  $p(\mathbf{x}_t | \mathbf{x}_T)$  satisfied is :

$$\begin{aligned} d\mathbf{x}_t &= [\mathbf{f}(\mathbf{x}_t, t) + g_t^2 \mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T)] dt + g_t d\mathbf{w}_t \\ &= \left( \theta_t + g_t^2 \frac{e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) (\mathbf{x}_T - \mathbf{x}_t) dt + g_t d\mathbf{w}_t \end{aligned} \quad (23)$$

Furthermore, we can derive the following transition probability of  $\mathbf{x}_t$  using Bayes' formula:

$$\begin{aligned} p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T) &= \frac{p(\mathbf{x}_T | \mathbf{x}_t, \mathbf{x}_0) p(\mathbf{x}_t | \mathbf{x}_0)}{p(\mathbf{x}_T | \mathbf{x}_0)} \\ &= \frac{p(\mathbf{x}_T | \mathbf{x}_t) p(\mathbf{x}_t | \mathbf{x}_0)}{p(\mathbf{x}_T | \mathbf{x}_0)} \end{aligned} \quad (24)$$

Since each component is independently and identically distributed (i.i.d), by considering a single dimension, we have:

$$\begin{aligned} p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T) &\propto \frac{1}{\sqrt{2\pi} \bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T} \exp - \left\{ \frac{(\mathbf{x}_t - [\mathbf{x}_T + (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_t}])^2}{2\bar{\sigma}_t^2} + \frac{(\mathbf{x}_T - [\mathbf{x}_T + (\mathbf{x}_t - \mathbf{x}_T) e^{-\bar{\theta}_{t:T}}])^2}{2\bar{\sigma}_{t:T}^2} \right\} \\ &= \frac{1}{\sqrt{2\pi} \bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T} \exp - \left\{ \frac{(\mathbf{x}_t - [\mathbf{x}_T + (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_t}])^2}{2\bar{\sigma}_t^2} + \frac{(\mathbf{x}_t - \mathbf{x}_T)^2 e^{-2\bar{\theta}_{t:T}}}{2\bar{\sigma}_{t:T}^2} \right\} \\ &\propto \frac{1}{\sqrt{2\pi} \bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T} \exp - \left\{ \left( \frac{1}{2\bar{\sigma}_t^2} + \frac{e^{-2\bar{\theta}_{t:T}}}{2\bar{\sigma}_{t:T}^2} \right) \mathbf{x}_t^2 - \left( \frac{\mathbf{x}_T - (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_t}}{\bar{\sigma}_t^2} + \frac{\mathbf{x}_T e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) \mathbf{x}_t \right\} \end{aligned} \quad (25)$$

Notice that:

$$\begin{aligned}
 \frac{1}{2\bar{\sigma}_t^2} + \frac{e^{-2\bar{\theta}_{t:T}}}{2\bar{\sigma}_{t:T}^2} &= \frac{\sigma_{t:T}^2 + \bar{\sigma}_t^2 e^{-2\bar{\theta}_{t:T}}}{2\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2} \\
 &= \frac{\lambda^2 \left[ (1 - e^{-2\bar{\theta}_{t:T}}) + (1 - e^{-2\bar{\theta}_t}) e^{-2\bar{\theta}_{t:T}} \right]}{2\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2} \\
 &= \frac{\lambda^2 \left[ (1 - e^{-2\bar{\theta}_{t:T}}) + (e^{-2\bar{\theta}_{t:T}} - e^{-2\bar{\theta}_T}) \right]}{2\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2} \\
 &= \frac{\bar{\sigma}_T^2}{2\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2}
 \end{aligned} \tag{26}$$

Bringing it back to (25), squaring the terms and reorganizing the equation, we obtain:

$$\begin{aligned}
 p(\mathbf{x}_t \mid \mathbf{x}_0, \mathbf{x}_T) &\propto \frac{1}{\sqrt{2\pi} \bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T} \exp - \left\{ \frac{\bar{\sigma}_T^2}{2\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2} \mathbf{x}_t^2 - \left( \frac{\mathbf{x}_T - (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_t}}{\bar{\sigma}_t^2} + \frac{\mathbf{x}_T e^{-2\bar{\theta}_{t:T}}}{\bar{\sigma}_{t:T}^2} \right) \mathbf{x}_t \right\} \\
 &= \frac{1}{\sqrt{2\pi} \bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T} \exp - \left\{ \frac{\mathbf{x}_t^2 - \left( \left[ \mathbf{x}_T - (\mathbf{x}_0 - \mathbf{x}_T) e^{-\bar{\theta}_t} \right] \frac{2\bar{\sigma}_{t:T}^2}{\bar{\sigma}_t^2} + e^{-2\bar{\theta}_{t:T}} \frac{2\bar{\sigma}_t^2}{\bar{\sigma}_T^2} \mathbf{x}_T \right) \mathbf{x}_t}{2(\bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T)^2} \right\} \\
 &\propto \frac{1}{\sqrt{2\pi} \bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T} \exp - \left\{ \frac{\mathbf{x}_t - e^{-\bar{\theta}_t} \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{x}_0 - \left[ \left( 1 - e^{-\bar{\theta}_t} \right) \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_t^2} + e^{-2\bar{\theta}_{t:T}} \frac{\bar{\sigma}_t^2}{\bar{\sigma}_T^2} \right] \mathbf{x}_T}{2(\bar{\sigma}_t \bar{\sigma}_{t:T} / \bar{\sigma}_T)^2} \right\}^2 \\
 &= N \left( e^{-\bar{\theta}_t} \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{x}_0 + \left[ \left( 1 - e^{-\bar{\theta}_t} \right) \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_t^2} + e^{-2\bar{\theta}_{t:T}} \frac{\bar{\sigma}_t^2}{\bar{\sigma}_T^2} \right] \mathbf{x}_T, \frac{\bar{\sigma}_t^2 \bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{I} \right)
 \end{aligned} \tag{27}$$

This concludes the proof of the **Proposition 3.1**.

## A.2. Proof of Proposition 3.3

**Proposition 3.3.** Let  $\mathbf{x}_t$  be a finite random variable describing by the given generalized Ornstein-Uhlenbeck process (4), for a fixed  $\mathbf{x}_T$ , the expectation of log-likelihood  $\mathbb{E}_{p(\mathbf{x}_0)}[\log p_\theta(\mathbf{x}_0 \mid \mathbf{x}_T)]$  possesses an Evidence Lower Bound (ELBO):

$$ELBO = \mathbb{E}_{p(\mathbf{x}_0)} \left[ \mathbb{E}_{p(\mathbf{x}_1 \mid \mathbf{x}_0)} [\log p_\theta(\mathbf{x}_0 \mid \mathbf{x}_1, \mathbf{x}_T)] - \sum_{t=2}^T KL(p(\mathbf{x}_{t-1} \mid \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T) \parallel p_\theta(\mathbf{x}_{t-1} \mid \mathbf{x}_t, \mathbf{x}_T)) \right] \tag{11}$$

Assuming  $p_\theta(\mathbf{x}_{t-1} \mid \mathbf{x}_t, \mathbf{x}_T)$  is a Gaussian distribution with a constant variance  $N(\boldsymbol{\mu}_{\theta, t-1}, \sigma_{\theta, t-1}^2 \mathbf{I})$ , maximizing the ELBO is equivalent to minimizing:

$$\mathcal{L} = \mathbb{E}_{t, \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T} \left[ \frac{1}{2\sigma_{\theta, t-1}^2} \|\boldsymbol{\mu}_{t-1} - \boldsymbol{\mu}_{\theta, t-1}\|^2 \right], \tag{12}$$

where  $\boldsymbol{\mu}_{t-1}$  represents the mean of  $p(\mathbf{x}_{t-1} \mid \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T)$ :

$$\boldsymbol{\mu}_{t-1} = \frac{1}{\bar{\sigma}_t'^2} [\bar{\sigma}_{t-1}'^2 (\mathbf{x}_t - b\mathbf{x}_T) a + (\bar{\sigma}_t'^2 - \bar{\sigma}_{t-1}'^2 a^2) \bar{\mathbf{m}}_t'], \tag{13}$$

where,

$$\begin{aligned}
 a &= \frac{e^{-\bar{\theta}_{t-1:t}} \bar{\sigma}_{t:T}^2}{\bar{\sigma}_{t-1:T}^2}, \\
 b &= \frac{1}{\bar{\sigma}_T^2} \left\{ (1 - e^{-\bar{\theta}_t}) \bar{\sigma}_{t:T}^2 + e^{-2\bar{\theta}_{t:T}} \bar{\sigma}_t^2 - \left[ (1 - e^{-\bar{\theta}_{t-1}}) \bar{\sigma}_{t-1:T}^2 + e^{-2\bar{\theta}_{t-1:T}} \bar{\sigma}_{t-1}^2 \right] a \right\}
 \end{aligned}$$

*Proof:* Firstly, followed by the theorem in DDPM (Ho et al., 2020):

$$\begin{aligned} \mathbb{E}_{p(\mathbf{x}_0)} [\log p_{\theta}(\mathbf{x}_0)] &\geq \mathbb{E}_{p(\mathbf{x}_0)} \left[ -KL(p(\mathbf{x}_T | \mathbf{x}_0) || p(\mathbf{x}_T)) + \mathbb{E}_{p(\mathbf{x}_1 | \mathbf{x}_0)} [\log p_{\theta}(\mathbf{x}_0 | \mathbf{x}_1)] \right. \\ &\quad \left. - \sum_{t=2}^T \mathbb{E}_{p(x_t | x_0)} [KL(p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t) || p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t))] \right] \end{aligned} \quad (28)$$

Similarly, we have:

$$\begin{aligned} \mathbb{E}_{p(\mathbf{x}_0)} [\log p_{\theta}(\mathbf{x}_0 | \mathbf{x}_T)] &\geq \mathbb{E}_{p(\mathbf{x}_0)} \left[ -KL(p(\mathbf{x}_T | \mathbf{x}_0, \mathbf{x}_T) || p(\mathbf{x}_T | \mathbf{x}_T)) + \mathbb{E}_{p(\mathbf{x}_1 | \mathbf{x}_0)} [\log p_{\theta}(\mathbf{x}_0 | \mathbf{x}_1, \mathbf{x}_T)] \right. \\ &\quad \left. - \sum_{t=2}^T \mathbb{E}_{p(x_t | x_0)} [KL(p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T) || p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{x}_T))] \right] \\ &= \mathbb{E}_{p(\mathbf{x}_0)} \left[ \mathbb{E}_{p(\mathbf{x}_1 | \mathbf{x}_0)} [\log p_{\theta}(\mathbf{x}_0 | \mathbf{x}_1, \mathbf{x}_T)] \right. \\ &\quad \left. - \sum_{t=2}^T \mathbb{E}_{p(x_t | x_0)} [KL(p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T) || p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{x}_T))] \right] \\ &= ELBO \end{aligned} \quad (29)$$

From Bayes' formula, we can infer that:

$$\begin{aligned} p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T) &= \frac{p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_{t-1}, \mathbf{x}_T)p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_T)}{p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)} \\ &= \frac{p(\mathbf{x}_t | \mathbf{x}_{t-1}, \mathbf{x}_T)p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_T)}{p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)} \end{aligned} \quad (30)$$

Since  $p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_T)$  and  $p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)$  are Gaussian distributions (8), by employing the reparameterization technique:

$$\begin{aligned} \mathbf{x}_{t-1} &= e^{-\bar{\theta}_{t-1}} \frac{\bar{\sigma}_{t-1:T}^2}{\bar{\sigma}_T^2} \mathbf{x}_0 + \left[ \left(1 - e^{-\bar{\theta}_{t-1}}\right) \frac{\bar{\sigma}_{t-1:T}^2}{\bar{\sigma}_T^2} + e^{-2\bar{\theta}_{t-1:T}} \frac{\bar{\sigma}_{t-1}^2}{\bar{\sigma}_T^2} \right] \mathbf{x}_T + \bar{\sigma}'_{t-1} \boldsymbol{\epsilon}_{t-1} \\ &= m(t-1)\mathbf{x}_0 + n(t-1)\mathbf{x}_T + \bar{\sigma}'_{t-1} \boldsymbol{\epsilon}_{t-1} \\ \mathbf{x}_t &= e^{-\bar{\theta}_t} \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} \mathbf{x}_0 + \left[ \left(1 - e^{-\bar{\theta}_t}\right) \frac{\bar{\sigma}_{t:T}^2}{\bar{\sigma}_T^2} + e^{-2\bar{\theta}_{t:T}} \frac{\bar{\sigma}_t^2}{\bar{\sigma}_T^2} \right] \mathbf{x}_T + \bar{\sigma}'_t \boldsymbol{\epsilon}_t \\ &= m(t)\mathbf{x}_0 + n(t)\mathbf{x}_T + \bar{\sigma}'_t \boldsymbol{\epsilon}_t \end{aligned} \quad (31)$$

Therefore,

$$\begin{aligned} \mathbf{x}_t &= \frac{m(t)}{m(t-1)} \mathbf{x}_{t-1} + \left[ n(t) - \frac{m(t)}{m(t-1)} n(t-1) \right] \mathbf{x}_T + \sqrt{\bar{\sigma}'_t{}^2 - \frac{m(t)^2}{m(t-1)^2} \bar{\sigma}'_{t-1}{}^2} \boldsymbol{\epsilon} \\ &= a\mathbf{x}_{t-1} + [n(t) - an(t-1)] \mathbf{x}_T + \sqrt{\bar{\sigma}'_t{}^2 - a^2 \bar{\sigma}'_{t-1}{}^2} \boldsymbol{\epsilon} \\ &= a\mathbf{x}_{t-1} + b\mathbf{x}_T + \sqrt{\bar{\sigma}'_t{}^2 - a^2 \bar{\sigma}'_{t-1}{}^2} \boldsymbol{\epsilon} \end{aligned} \quad (32)$$

Thus,  $p(\mathbf{x}_t | \mathbf{x}_{t-1}, \mathbf{x}_T) = N(a\mathbf{x}_{t-1} + b\mathbf{x}_T, (\bar{\sigma}'_t{}^2 - a^2 \bar{\sigma}'_{t-1}{}^2) \mathbf{I})$  is also a Gaussian distribution. Bring it back to equation (30) we can easily obtain :

$$\boldsymbol{\mu}_{t-1} = \frac{1}{\bar{\sigma}'_t{}^2} [\bar{\sigma}'_{t-1}{}^2 (\mathbf{x}_t - b\mathbf{x}_T) a + (\bar{\sigma}'_t{}^2 - \bar{\sigma}'_{t-1}{}^2 a^2) \bar{\mathbf{m}}'_t], \quad (13)$$

Accordingly,

$$\begin{aligned}
 & KL(p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T) || p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{x}_T)) \\
 &= \mathbb{E}_{p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T)} \left[ \log \frac{\frac{1}{\sqrt{2\pi}\sigma_{t-1}} e^{-(\mathbf{x}_{t-1} - \boldsymbol{\mu}_{t-1})^2 / 2\sigma_{t-1}^2}}{\frac{1}{\sqrt{2\pi}\sigma_{\theta,t-1}} e^{-(\mathbf{x}_{t-1} - \boldsymbol{\mu}_{\theta,t-1})^2 / 2\sigma_{\theta,t-1}^2}} \right] \\
 &= \mathbb{E}_{p(\mathbf{x}_{t-1} | \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T)} \left[ \log \sigma_{\theta,t-1} - \log \sigma_{t-1} - (\mathbf{x}_{t-1} - \boldsymbol{\mu}_{t-1})^2 / 2\sigma_{t-1}^2 + (\mathbf{x}_{t-1} - \boldsymbol{\mu}_{\theta,t-1})^2 / 2\sigma_{\theta,t-1}^2 \right] \\
 &= \log \sigma_{\theta,t-1} - \log \sigma_{t-1} - \frac{1}{2} + \frac{\sigma_{t-1}^2}{2\sigma_{\theta,t-1}^2} + \frac{(\boldsymbol{\mu}_{t-1} - \boldsymbol{\mu}_{\theta,t-1})^2}{2\sigma_{\theta,t-1}^2}
 \end{aligned} \tag{33}$$

Ignoring unlearnable constant, the training object that involves minimizing the negative ELBO is :

$$\mathcal{L} = \mathbb{E}_{t, \mathbf{x}_0, \mathbf{x}_t, \mathbf{x}_T} \left[ \frac{1}{2\sigma_{\theta,t-1}^2} \|\boldsymbol{\mu}_{t-1} - \boldsymbol{\mu}_{\theta,t-1}\|^2 \right], \tag{34}$$

This concludes the proof of the **Proposition 3.3**.

## B. Theoretical Results

### B.1. Brownian Bridge

In this section, we will show the mathematical essence of some other bridge models, some of which are all equivalent.

**Proposition B.1.** *The mathematical essence of BBDM (Li et al., 2023), DDBM (VE) (Zhou et al., 2024) and I<sup>2</sup>SB (Liu et al., 2023a) are all equivalent to the Brownian bridge.*

*Proof:* Firstly, it is easy to understand that BBDM uses the Brownian bridge as its fundamental model architecture.

The DDBM (VE) model is derived as the Doob's  $h$ -transform of VE-SDE, and we begin by specifying the SDE:

$$d\mathbf{x}_t = d\mathbf{w}_t \tag{35}$$

Its transition probability is given by:

$$p(\mathbf{x}_t | \mathbf{x}_s) = N(\mathbf{x}_s, t - s) \tag{36}$$

Since, the  $h$ -function of SDE (35) is:

$$\begin{aligned}
 \mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T) &= \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_T | \mathbf{x}_t) \\
 &= \frac{\mathbf{x}_T - \mathbf{x}_t}{T - t}
 \end{aligned} \tag{37}$$

Therefore, the Doob's  $h$ -transform of (35) is:

$$d\mathbf{x}_t = \frac{\mathbf{x}_T - \mathbf{x}_t}{T - t} dt + d\mathbf{w}_t \tag{38}$$

That is the definition of Brownian bridge. Hence, DDBM (VE) is a Brownian bridge model.

Furthermore, the transition kernel of (38) is:

$$\begin{aligned}
 p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T) &= \frac{p(\mathbf{x}_T | \mathbf{x}_t, \mathbf{x}_0)p(\mathbf{x}_t | \mathbf{x}_0)}{p(\mathbf{x}_T | \mathbf{x}_0)} \\
 &= \frac{p(\mathbf{x}_T | \mathbf{x}_t)p(\mathbf{x}_t | \mathbf{x}_0)}{p(\mathbf{x}_T | \mathbf{x}_0)} \\
 &= \frac{N(\mathbf{x}_t, T - t)N(\mathbf{x}_0, t)}{N(\mathbf{x}_0, T)} \\
 &= N\left(\left(1 - \frac{t}{T}\right)\mathbf{x}_0 + \frac{t}{T}\mathbf{x}_T, \frac{t(T-t)}{T}\mathbf{I}\right)
 \end{aligned} \tag{39}$$

This precisely corresponds to the sampling process of I<sup>2</sup>SB, thus confirming that I<sup>2</sup>SB also represents a Brownian bridge.

This concludes the proof of the **Proposition B.1**.

## B.2. Connections Between GOU, VE and VP

The following proposition will show us that both VE and VP processes are special cases of GOU process:

**Proposition 5.1.** *For a given GOU process (4), there exists relationships:*

$$\begin{aligned} \lim_{\theta_t \rightarrow 0} \text{GOU} &= \text{VE} \\ \lim_{\mu \rightarrow 0, \lambda \rightarrow 1} \text{GOU} &= \text{VP} \end{aligned} \quad (18)$$

*Proof:* It's easy to know:

$$\begin{aligned} \lim_{\theta_t \rightarrow 0} \text{GOU} &= \lim_{\theta_t \rightarrow 0} \{d\mathbf{x}_t = \theta_t (\boldsymbol{\mu} - \mathbf{x}_t) dt + g_t d\mathbf{w}_t\} \\ &= \lim_{\theta_t \rightarrow 0} \{d\mathbf{x}_t = g_t d\mathbf{w}_t\} \\ &= \text{VE}, \end{aligned} \quad (40)$$

where  $g_t$  will be controlled by  $\lambda^2$ .

Besides, we have:

$$\begin{aligned} \lim_{\mu \rightarrow 0, \lambda \rightarrow 1} \text{GOU} &= \lim_{\mu \rightarrow 0, \lambda \rightarrow 1} \{d\mathbf{x}_t = \theta_t (\boldsymbol{\mu} - \mathbf{x}_t) dt + g_t d\mathbf{w}_t\} \\ &= \lim_{\mu \rightarrow 0, \lambda \rightarrow 1} \{d\mathbf{x}_t = \theta_t \boldsymbol{\mu} dt - \theta_t \mathbf{x}_t dt + g_t d\mathbf{w}_t\} \\ &= \lim_{\mu \rightarrow 0, \lambda \rightarrow 1} \left\{ d\mathbf{x}_t = -\frac{1}{2} g_t^2 \mathbf{x}_t dt + g_t d\mathbf{w}_t \right\} \\ &= \text{VP}, \end{aligned} \quad (41)$$

where  $g_t$  will be controlled by  $\theta_t$ .

This concludes the proof of the **Proposition 5.1**.

## C. GOU Process

**Theorem C.1.** *For a given GOU process:*

$$d\mathbf{x}_t = \theta_t (\boldsymbol{\mu} - \mathbf{x}_t) dt + g_t d\mathbf{w}_t \quad (4)$$

where  $\boldsymbol{\mu}$  is a given state vector,  $\theta_t$  denotes a scalar drift coefficient and  $g_t$  represents the diffusion coefficient. It possesses a closed-form analytical solution:

$$p(\mathbf{x}_t | \mathbf{x}_s) = N \left( \boldsymbol{\mu} + (\mathbf{x}_s - \boldsymbol{\mu}) e^{-\bar{\theta}_{s:t}}, \frac{g_t^2}{2\theta_t} (1 - e^{-2\bar{\theta}_{s:t}}) \mathbf{I} \right), \quad \bar{\theta}_{s:t} = \int_s^t \theta_z dz \quad (5)$$

*Proof:* Writing:

$$\mathbf{f}(\mathbf{x}_t, t) = \mathbf{x}_t e^{\bar{\theta}_t} \quad (42)$$

Using Ito differential formula, we get:

$$\begin{aligned} d\mathbf{f}(\mathbf{x}_t, t) &= \mathbf{x}_t \theta_t e^{\bar{\theta}_t} dt + e^{\bar{\theta}_t} d\mathbf{x}_t \\ &= \mathbf{x}_t \theta_t e^{\bar{\theta}_t} dt + e^{\bar{\theta}_t} [\theta_t (\boldsymbol{\mu} - \mathbf{x}_t) dt + g_t d\mathbf{w}_t] \\ &= e^{\bar{\theta}_t} \theta_t \boldsymbol{\mu} + e^{\bar{\theta}_t} g_t d\mathbf{w}_t \end{aligned} \quad (43)$$

Integrating from  $s$  to  $t$  we get:

$$\begin{aligned} \mathbf{x}_t e^{\bar{\theta}_t} - \mathbf{x}_s e^{\bar{\theta}_s} &= \int_s^t e^{\bar{\theta}_z} \theta_z \boldsymbol{\mu} dz + \int_s^t e^{\bar{\theta}_z} g_z d\mathbf{w}_z \\ &= (e^{\bar{\theta}_t} - e^{\bar{\theta}_s}) \boldsymbol{\mu} + \int_s^t e^{\bar{\theta}_z} g_z d\mathbf{w}_z \end{aligned} \quad (44)$$



It's obvious that the transition kernel is a Gaussian distribution. Since  $d\mathbf{w}_z \sim N(\mathbf{0}, dz\mathbf{I})$ , we have:

$$\begin{aligned} \int_s^t e^{\bar{\theta}_z} g_z d\mathbf{w}_z &= N\left(\mathbf{0}, \int_s^t e^{2\bar{\theta}_z} g_z^2 dz \mathbf{I}\right) \\ &= N\left(\mathbf{0}, \lambda^2 \int_s^t e^{2\bar{\theta}_z} 2\theta_t dz \mathbf{I}\right) \\ &= N\left(\mathbf{0}, \lambda^2 \left(e^{2\bar{\theta}_t} - e^{2\bar{\theta}_s}\right) \mathbf{I}\right) \end{aligned} \quad (45)$$

Therefore:

$$\begin{aligned} \mathbf{x}_t e^{\bar{\theta}_t} - \mathbf{x}_s e^{\bar{\theta}_s} &= \left(e^{\bar{\theta}_t} - e^{\bar{\theta}_s}\right) \boldsymbol{\mu} + N\left(\mathbf{0}, \lambda^2 \left(e^{2\bar{\theta}_t} - e^{2\bar{\theta}_s}\right) \mathbf{I}\right) \\ \mathbf{x}_t &= \boldsymbol{\mu} + (\mathbf{x}_s - \boldsymbol{\mu}) e^{-\bar{\theta}_{s:t}} + N\left(\mathbf{0}, \frac{g_t^2}{2\theta_t} \left(1 - e^{-2\bar{\theta}_{s:t}}\right) \mathbf{I}\right) \end{aligned} \quad (46)$$

This concludes the proof of the **Theorem C.1**.

## D. Doob's $h$ -transform

**Theorem D.1.** For a given SDE:

$$d\mathbf{x}_t = \mathbf{f}(\mathbf{x}_t, t) dt + g_t d\mathbf{w}_t, \quad \mathbf{x}_0 \sim p(\mathbf{x}_0), \quad (1)$$

For a fixed  $\mathbf{x}_T$ , the evolution of conditional probability  $p(\mathbf{x}_t | \mathbf{x}_T)$  follows:

$$d\mathbf{x}_t = [\mathbf{f}(\mathbf{x}_t, t) + g_t^2 \mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T)] dt + g_t d\mathbf{w}_t, \quad \mathbf{x}_0 \sim p(\mathbf{x}_0 | \mathbf{x}_T), \quad (6)$$

where  $\mathbf{h}(\mathbf{x}_t, t, \mathbf{x}_T, T) = \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_T | \mathbf{x}_t)$ .

*Proof:*  $p(\mathbf{x}_t | \mathbf{x}_0)$  satisfies Kolmogorov Forward Equation (KFE) also called Fokker-Planck equation (Risken & Risken, 1996):

$$\frac{\partial}{\partial t} p(\mathbf{x}_t | \mathbf{x}_0) = -\nabla_{\mathbf{x}_t} \cdot [\mathbf{f}(\mathbf{x}_t, t) p(\mathbf{x}_t | \mathbf{x}_0)] + \frac{1}{2} g_t^2 \nabla_{\mathbf{x}_t} \cdot \nabla_{\mathbf{x}_t} p(\mathbf{x}_t | \mathbf{x}_0) \quad (47)$$

Similarly,  $p(\mathbf{x}_T | \mathbf{x}_t)$  satisfies Kolmogorov Backward Equation (KBE) (Risken & Risken, 1996):

$$-\frac{\partial}{\partial t} p(\mathbf{x}_T | \mathbf{x}_t) = \mathbf{f}(\mathbf{x}_t, t) \cdot \nabla_{\mathbf{x}_t} p(\mathbf{x}_T | \mathbf{x}_t) + \frac{1}{2} g_t^2 \nabla_{\mathbf{x}_t} \cdot \nabla_{\mathbf{x}_t} p(\mathbf{x}_T | \mathbf{x}_t) \quad (48)$$

Using Bayes' rule, we have:

$$\begin{aligned} p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T) &= \frac{p(\mathbf{x}_T | \mathbf{x}_t, \mathbf{x}_0) p(\mathbf{x}_t | \mathbf{x}_0)}{p(\mathbf{x}_T | \mathbf{x}_0)} \\ &= \frac{p(\mathbf{x}_T | \mathbf{x}_t) p(\mathbf{x}_t | \mathbf{x}_0)}{p(\mathbf{x}_T | \mathbf{x}_0)} \end{aligned} \quad (49)$$



This is the definition of FP equation of conditional transition probability  $p(\mathbf{x}_t | \mathbf{x}_0, \mathbf{x}_T)$ , which represents the evolution follows the SDE:

$$d\mathbf{x}_t = [\mathbf{f}(\mathbf{x}_t, t) + g_t^2 \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_T | \mathbf{x}_t)] dt + g_t d\mathbf{w}_t \quad (53)$$

This concludes the proof of the **Theorem D.1**.

## E. Experimental Details

For all experiments, we use the same noise network, with the network architecture and mainly training parameters consistent with the paper (Luo et al., 2023a). This network is similar to a U-Net structure but without group normalization layers and self-attention layers. The steady variance level  $\lambda^2$  was set to 30 (over 255), and the sampling step number T was set to 100. In the training process, we set the patch size = 128 with batch size = 8 and use Adam (Kingma & Ba, 2015) optimizer with parameters  $\beta_1 = 0.9$  and  $\beta_2 = 0.99$ . The total training steps are 900 thousand with the initial learning rate set to  $10^{-4}$ , and it decays by half at iterations 300, 500, 600, and 700 thousand. For the setting of  $\theta_t$ , we employ a flipped version of cosine noise schedule (Nichol & Dhariwal, 2021), enabling  $\theta_t$  to change from 0 to 1 over time. Notably, to address the issue of  $\theta_t$  being too smooth when  $t$  closed to 1, we let the coefficient  $e^{-\theta_T}$  to be a small enough value  $\delta = 0.005$  instead of zero, which represents  $\bar{\theta}_T \approx \sum_{i=0}^T \theta_i dt = -\log \delta$ , as well as  $dt = -\log \delta / \sum_{i=0}^T \theta_i$ . Our models are trained on a single 3090 GPU with 24GB memory for about 2.5 days.

## F. Additional Experiments

Table 6. **Image Inpainting**. Qualitative comparison with the relevant baselines on CelebA-HQ with thick mask.

METHOD	PSNR $\uparrow$	SSIM $\uparrow$	LPIPS $\downarrow$	FID $\downarrow$
DDRM	19.48	0.8154	0.1487	26.24
IRSDE	21.12	0.8499	0.1046	11.12
GOUB	<b>22.27</b>	<b>0.8754</b>	<b>0.0914</b>	<b>5.64</b>

Table 7. **Image Deraining**. Qualitative comparison with the relevant baselines on Rain100L.

METHOD	PSNR $\uparrow$	SSIM $\uparrow$	LPIPS $\downarrow$	FID $\downarrow$
PRENET	37.48	0.9792	0.020	10.9
MAXIM	38.06	0.9770	0.048	19.0
IRSDE	38.30	0.9805	0.014	7.94
GOUB	<b>39.79</b>	<b>0.9830</b>	<b>0.009</b>	<b>5.18</b>

Table 8. **Image 8 $\times$  Super-Resolution**. Qualitative comparison with the relevant baselines on DIV2K.

METHOD	PSNR $\uparrow$	SSIM $\uparrow$	LPIPS $\downarrow$	Training Datasets
SRFlow	23.05	0.57	<b>0.272</b>	DIV2K + Flickr2K
IRSDE	22.34	0.55	0.331	DIV2K
GOUB	<b>23.17</b>	<b>0.60</b>	0.310	DIV2K

## G. Additional Visual Results

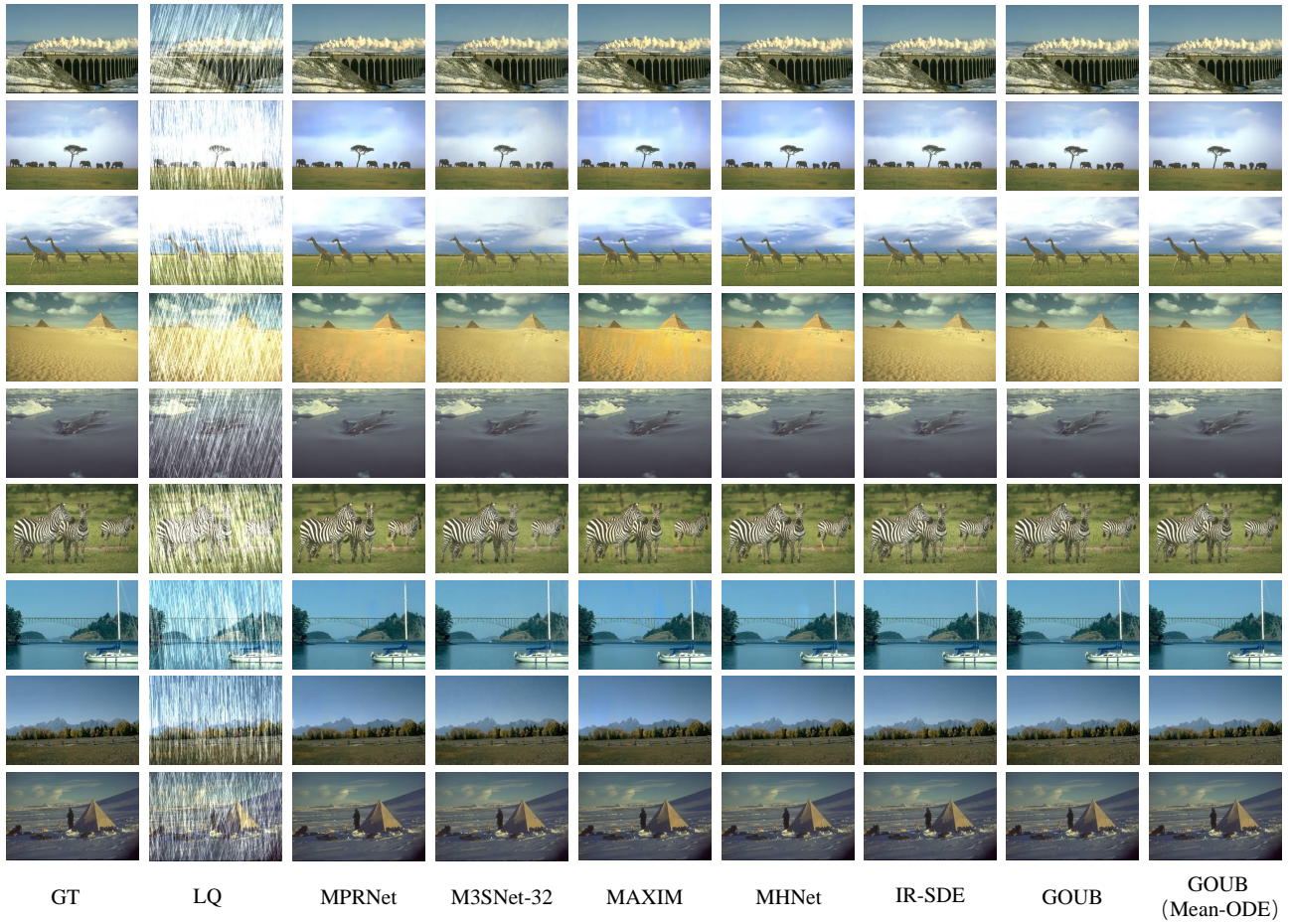


Figure 6. Additional visual results on deraining with Rain100H datasets.

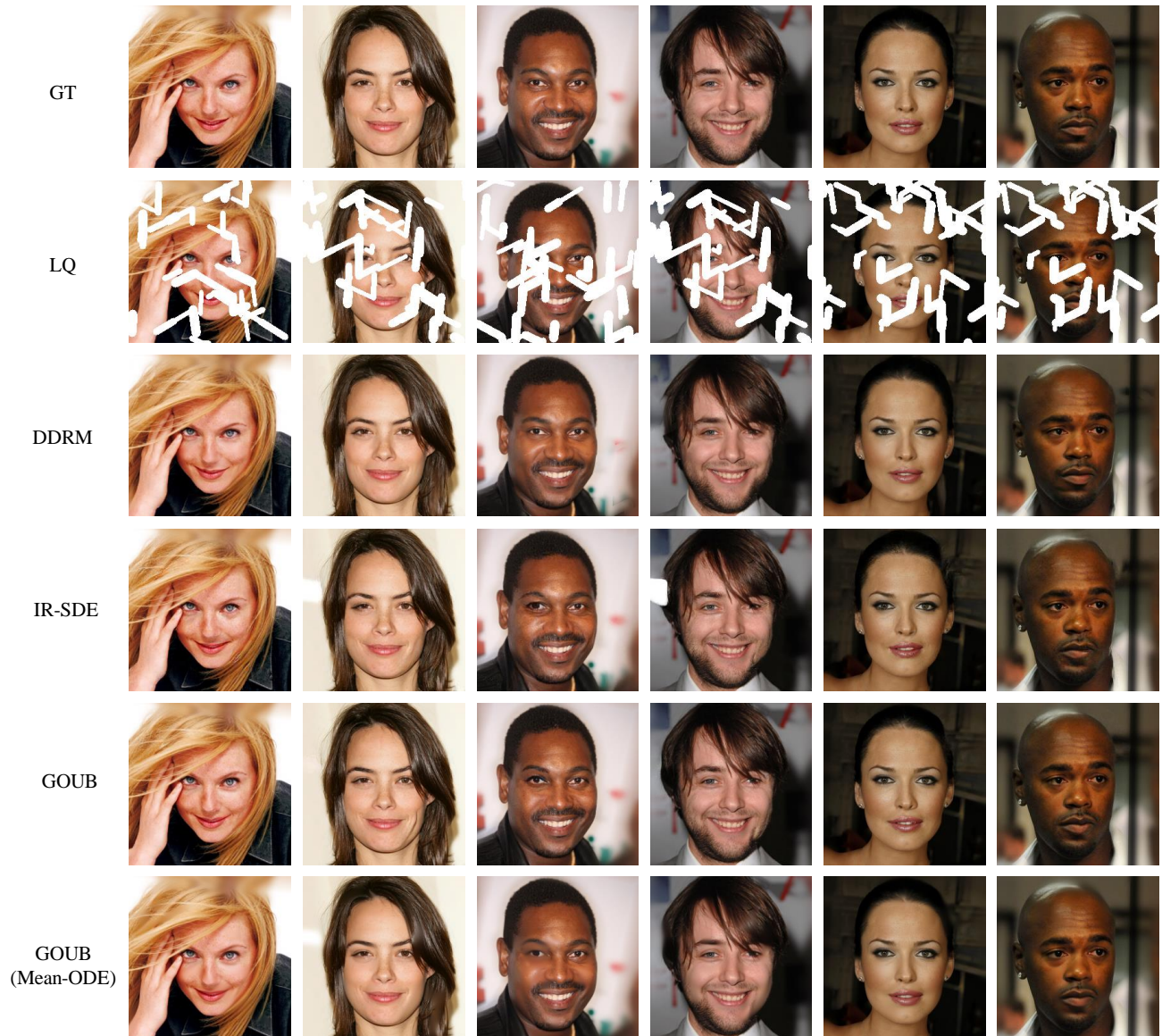


Figure 7. Additional visual results on thin mask inpainting with CelebA-HQ datasets.

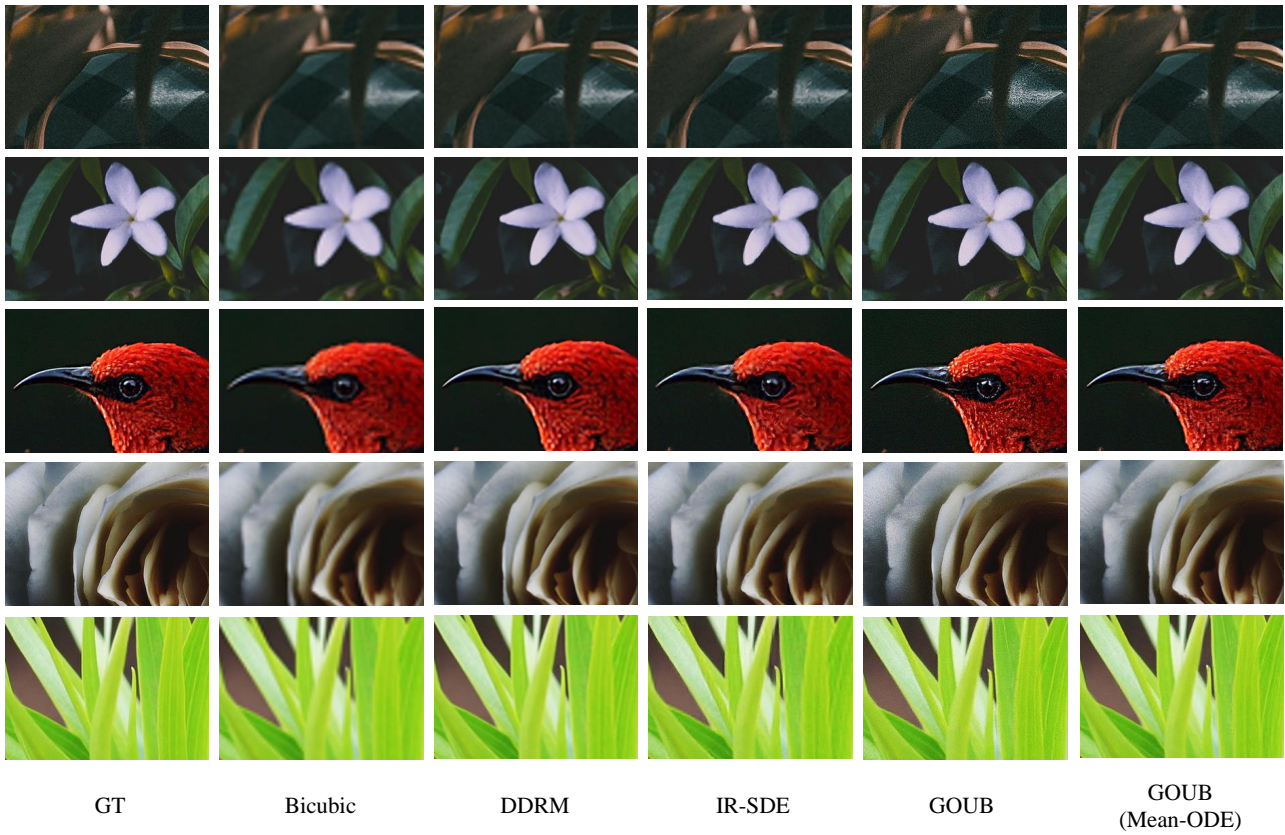


Figure 8. Additional visual results on 4x super-resolution with DIV2K datasets.