000 **POGDIFF: PRODUCT-OF-GAUSSIANS DIFFUSION MOD-**001 ELS FOR IMBALANCED TEXT-TO-IMAGE GENERATION 002 003

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ABSTRACT

Diffusion models have made significant advancements in recent years. However, their performance often deteriorates when trained or fine-tuned on imbalanced datasets. This degradation is largely due to the disproportionate representation of majority and minority data in image-text pairs. In this paper, we propose a general fine-tuning approach, dubbed PoGDiff, to address this challenge. Rather than directly minimizing the KL divergence between the predicted and ground-truth distributions, PoGDiff replaces the ground-truth distribution with a Product of Gaussians (PoG), which is constructed by combining the original ground-truth targets with the predicted distribution conditioned on a neighboring text embedding. Experiments on real-world datasets demonstrate that our method effectively addresses the imbalance problem in diffusion models, improving both generation accuracy and quality.

- 1 INTRODUCTION
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028 The development of diffusion models (Ho et al., 2020; Song et al., 2020b) and their subsequent exten-029 sions (Song et al., 2020a; Nichol & Dhariwal, 2021; Huang et al., 2023) has significantly advanced the learning of complex probability distributions across various data types, including images (Ho et al., 031 2022; Rombach et al., 2022; Saharia et al., 2022; Ho & Salimans, 2022), audio (Kong et al., 2020), and 3D biomedical imaging data (Luo & Hu, 2021; Poole et al., 2022; Shi et al., 2023; Pinaya et al., 2022). For these generative models, the amount of training data plays a critical role in determining 033 both the accuracy of probability estimation and the model's ability to generalize, which enables 034 effective extrapolation within the probability space. 035

Data diversity and abundance are key to improving the generative capabilities of large-scale models, 037 enabling them to capture intricate details within a vast probability space. However, many data-driven 038 modeling tasks often rely on small, imbalanced real-world datasets, leading to poor generation quality, particularly for minority groups. For example, when training and fine-tuning a diffusion model with an imbalanced dataset of individuals, existing models often struggle to generate accurate images for 040 those who appear less frequently (i.e., minorities) in the training data (Fig. 1). 041

042 This limitation is true even for finetuning large diffusion models pretrained on large-scale datasets 043 like LAION-5B (Schuhmann et al., 2022), e.g., Stable Diffusion (Rombach et al., 2022). Imagine 044 an imbalanced dataset consisting of employees in a small company, senior employees might have more photos available, while new employees only have a very limited number of them. Since none of the employees appear in the LAION-5B dataset, generating photos of them require finetuning the 046 Stable Diffusion model. Unfortunately, finetuning the model on such an imbalanced dataset might 047 enable the model to generate accurate images for the majority group (i.e., senior employees), but it 048 will perform poorly for the minority group (i.e., new employees). 049

To address this challenge, we propose a general fine-tuning approach, dubbed PoGDiff. Rather than 051 directly minimizing the KL divergence between the predicted and ground-truth distributions, PoGDiff replaces the ground-truth distribution with a Product of Gaussians (PoG), which is constructed 052 by combining the original ground-truth targets with the predicted distribution conditioned on a neighboring text embedding. Our contributions are as follows:

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Figure 1: **PoGDiff for imbalanced text-to-image (IT2I) generation.** Existing methods such as Stable Diffusion (Rombach et al., 2022) and CBDM (Qin et al., 2023) fall short for minority data (**Low Density**). In contrast, Our PoGDiff successfully generates high-quality images even for minority data, outperforming all the baselines.

- We identify the problem of imbalanced text-to-image generation (IT2I) and introduce the first general diffusion model, dubbed Product-of-Gaussians Diffusion Models (PoGDiff), for addressing this problem.
- Our theoretical analysis shows that training of PoGDiff is equivalent to training a normal diffusion model while encouraging the model to generate the same image given similar text prompts (conditions).
- Our empirical results on real-world datasets demonstrate the effectiveness of our method, outperforming all state-of-the-art baselines.

2 RELATED WORK

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Long-Tailed Recognition. Addressing the challenges posed by long-tailed data distributions has been a critical area of research in machine learning, for both classification and regression problems. 087 Traditional methods, such as re-sampling and re-weighting techniques, have been used to mitigate 880 class imbalances by either over-sampling minority classes or assigning higher weights to them during 089 training (Chawla et al., 2002; He & Garcia, 2009; Torgo et al., 2013; Branco et al., 2017; 2018). 090 Such algorithms fail to measure the distance in continuous label space and fall short in handling 091 high-dimensional data (e.g., images and text). Deep imbalanced regression methods (Yang et al., 092 2021; Ren et al., 2022; Gong et al., 2022; Keramati et al., 2023; Wang & Wang, 2024) address this challenge by reweighting the data using the effective label density during representation learning. 093 However, all these methods above are designed for *recognition* tasks such as classification and 094 regression, and are therefore not applicable to our generation task. 095

Diffusion Models Related to Long-Tailed Data. There are also works that related to both diffusion models and long-tailed data. They aim at improving generation robustness using noisy label (Na et al., 2024), improving fairness in image generation (Shen et al., 2023), and improving classification accuracy using diffusion models (Zhang et al., 2024). However, these works have different goals and therefore are not applicable to our setting.

Most relevant to our work is Class Balancing Diffusion Model (CBDM) (Qin et al., 2023), which uses a distribution adjustment regularizer that enhances tail-class generation based on the model's predictions for the head class. It improves the quality of long-tailed generation by assuming one-hot conditional labels (i.e., classification-based settings). However, this assumption does not generalize to the modern setting where image generation is usually conditioned on free-form text prompts. As a result, when adapted to the free-form setting, they often fail to model the similarity among different text prompts, leading to suboptimal generation performance in minority data (as verified by empirical results in Sec. 4).

108 3 **METHODS**

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110 3.1 Preliminaries 111

112 Diffusion models (DMs) (Ho et al., 2020) are probabilistic models that generate an output image \mathbf{x}_0 from a random noise vector \mathbf{x}_T conditioned on text input c. DMs operate through two main 113 processes: the forward diffusion process and the reverse denoising process. During the diffusion 114 process, Gaussian noise is progressively added to a data sample x_0 over T steps. The forward process 115 is defined as a Markov chain, where: 116

$$q\left(\mathbf{x}_{t}|\mathbf{x}_{t-1}\right) = \mathcal{N}\left(\mathbf{x}_{t}; \sqrt{1-\beta_{t}}\mathbf{x}_{t-1}, \beta_{t}\mathbf{I}\right)$$

119 Here, β_t is the predefined diffusion rate at step t. By denoting $\alpha_t = 1 - \beta_t$ and $\bar{\alpha}_t = \prod_{i=1}^t \alpha_i$, we 120 can describe the entire diffusion process as: 121

$$q\left(\mathbf{x}_{1:T}|\mathbf{x}_{0}\right) = \prod_{t=1}^{T} q\left(\mathbf{x}_{t}|\mathbf{x}_{t-1}\right)$$
$$q\left(\mathbf{x}_{t}|\mathbf{x}_{0}\right) = \mathcal{N}\left(\mathbf{x}_{t}; \sqrt{\bar{\alpha}_{t}}\mathbf{x}_{0}, (1-\bar{\alpha}_{t})\mathbf{I}\right)$$

125 The denoising process removes noise from the sample x_T , eventually recovering x_0 . A denoising 126 model $\epsilon_{\theta}(\mathbf{x}_t, t, \mathbf{y})$ is trained to estimate the noise ϵ from \mathbf{x}_t and a text-guided embedding $\mathbf{y} = \phi(\mathbf{c})$, 127 where $\phi(\cdot)$ is a pretrained text encoder. Formally: 128

 $p_{\theta}\left(\mathbf{x}_{t-1} | \mathbf{x}_{t}, t, \mathbf{y}\right) = \mathcal{N}\left(\mathbf{x}_{t-1}; \epsilon_{\theta}(\mathbf{x}_{t}, t, \mathbf{y}), \sigma_{t}^{2} \mathbf{I}\right).$

130 The denoising process is trained by maximizing the likelihood of the data under the model or, 131 equivalently, by minimizing the variational lower bound on the negative log-likelihood of the data. Ho et al. (2020) shows that this is equivalent to minimizing the KL divergence between the predicted 132 distribution $p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{y})$ and the ground-truth distribution $q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0, \mathbf{y})$ at each time step t 133 during the backward process. The training objective then becomes: 134

$$\min D_{KL}\left(q\left(\mathbf{x}_{t-1} | \mathbf{x}_{t}, \mathbf{x}_{0}, \mathbf{y}\right) \| p_{\theta}\left(\mathbf{x}_{t-1} | \mathbf{x}_{t}, \mathbf{y}\right)\right).$$

This can be simplified to:

$$L_{DM} = \mathbb{E}_{\mathbf{x}_0 = \mathbf{x}, \epsilon \sim \mathcal{N}(0, \mathbf{I}), t} \left[\|\epsilon - \epsilon_{\theta}(\mathbf{x}_t, t, \mathbf{y})\|_2^2 \right]$$

Latent diffusion models (LDMs) (Rombach et al., 2022) are diffusion models that perform the entire 140 diffusion and denoising process in a lower-dimensional latent space. LDMs first learn an encoder 141 \mathcal{E} and a decoder \mathcal{D} , which are then frozen during subsequent training of the diffusion models. The 142 corresponding objective is then simplified to: 143

$$L_{LDM} = \mathbb{E}_{\mathbf{z}_0 = \mathcal{E}(\mathbf{x}), \epsilon \sim \mathcal{N}(0, \mathbf{I}), t} \left[\|\epsilon - \epsilon_{\theta}(\mathbf{z}_t, t, \mathbf{y})\|_2^2 \right]$$

In this paper, we use Stable Diffusion (LDM) (Rombach et al., 2022) as our backbone model. Since 146 our method works for both the vanilla DMs and LDMs, for clarity, we use the notation \mathbf{x} instead of \mathbf{z} , as the encoder \mathcal{E} and decoder \mathcal{D} are fixed during fine-tuning. 148

149 3.2 PRODUCT-OF-GAUSSIANS DIFFUSION MODELS (POGDIFF) 150

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       3.2.1 MAIN IDEA
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152 Method Overview. Given an image dataset 153 $\mathcal{D} = {\mathbf{x}^{(i)}, \mathbf{c}^{(i)}}_{i=1}^{N}$, where $\mathbf{c}^{(i)}$ is the text description for image $\mathbf{x}^{(i)}$, we use a fixed CLIP 154 155 encoder to produce $c^{(i)}$'s corresponding text em-156 bedding $\mathbf{y} = \phi(\mathbf{c})$. 157

158 Typical diffusion models minimize the KL 159 divergence between the predicted distribution $p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{y}) = \mathcal{N}(\epsilon_{\theta}(\mathbf{x}_t, t, \mathbf{y}), \lambda_{\mathbf{y}}^{-1}\mathbf{I})$ 160 and the ground-truth distribution 161 $q(\mathbf{x}_{t-1}|\mathbf{x}_t,\mathbf{x}_0,\mathbf{y}) = \mathcal{N}(\epsilon,\lambda_t^{-1}\mathbf{I})$ at each



Figure 2: Overview of our PoGDiff. During finetuning, PoGDiff collects k neighbors of the current text embedding y and samples one y' from them based on Eqn. (8). Both y and y' will then be employed to denoise the current image x_t to x_{t-1} .

time step t during the backward denoising process. Here, λ_{y} and λ_{t} represent the precision. In contrast, our PoGDiff replaces the ground-truth target with a Product of Gaussians (PoG), and instead minimize the following KL divergence (for each t)

$$\mathcal{L}_{t-1}^{\text{PoGDiff}} = D_{KL} \left(q \left(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{x}_0, \mathbf{y} \right) \circ p_{\theta} \left(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{y}' \right) \left\| p_{\theta} \left(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{y} \right) \right), \tag{1}$$

where \circ represents the product of two Gaussian distributions, \mathbf{y}' is a selected neighboring embedding from other samples in the training dataset (more details below), and $p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{y}')$ denotes the predicted distribution when using \mathbf{y}' as the input text embedding.

170 As shown in Fig. 2, intuitively, PoGDiff's denoising model $\epsilon_{\theta}(\mathbf{x}_t, t, \mathbf{y})$ (or $p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{y})$) 171 is optimized towards two target distributions, 172 equivalently increasing the weights for minority 173 instances (more details below). This approach 174 enhances the text-to-image mapping by leverag-175 ing the statistical strength of neighboring data 176 points, thereby improving and quality of the gen-177 erated images, especially for minority images. 178

Intuition behind the Product of Gaussians 179 (PoG). During fine-tuning, typical diffusion 180 models "lock" the text conditional embedding 181 $\mathbf{y} = \psi(c)$ to the corresponding image x. Conse-182 quently, if the dataset follows a long-tailed dis-183 tribution, the fine-tuned or post-trained diffusion model becomes heavily biased toward the ma-185 jority data, performing poorly on minority data. 186 Fig. 3 demonstrates our intuition. When training using a text-image pair (\mathbf{y}, \mathbf{x}) , our PoGDiff "bor-187 rows" information from neighboring text condi-188 tional embedding y', thereby effectively increas-189 ing the data density in the minority region and 190 leading to smoother (less imbalanced) effective 191 density, as shown in Fig. 3 (right). However, 192 since the text embedding is fixed during fine-193 tuning (i.e., ϕ is frozen), directly smoothing the 194 text embedding space is not feasible. Instead, we 195 rely on the properties of the product of Gaussian 196 distributions.



Figure 3: Comparing denoising networks of typical diffusion models (Ho et al., 2020; Rombach et al., 2022) and our PoGDiff. Left: In conditional text-to-image diffusion models, a data point (i.e., \mathbf{x}) is mainly affected by its text embedding (also affected by the random latent codes). Right: In PoGDiff, neighbors participate to modulate the final effective density. Here, \mathbf{y} denotes the text prompts, which are the embeddings of the text descriptions of the images; \mathbf{x} denotes the associated images. The tightly packed circles at the top indicate higher density, while the sparse circles indicate lower density.

By definition, given two Gaussian distributions, $\mathcal{N}(\mu_1, \lambda_1^{-1}\mathbf{I})$ and $\mathcal{N}(\mu_2, \lambda_2^{-1}\mathbf{I})$, their product is still a Gaussian distribution:

$$\mathcal{N}(\mu_1, \lambda_1^{-1}) \circ \mathcal{N}(\mu_2, \lambda_2^{-1}) = \mathcal{N}\left(\frac{\lambda_1 \mu_1 + \lambda_2 \mu_2}{\lambda_1 + \lambda_2}, (\lambda_1 + \lambda_2)^{-1}\right) \triangleq \mathcal{N}\left(\mu_{\text{PoG}}, \lambda_{\text{PoG}}^{-1}\right), \quad (2)$$

which can be treated as a "composition" of two individual Gaussians, incorporating information from both. This intuition is key to developing our PoGDiff objective function.

3.2.2 THEORETICAL ANALYSIS AND ALGORITHMIC DESIGN

Based on Eqn. (1), we then derive a concrete objective function following Proposition 3.1 below.

Proposition 3.1. Assume $\lambda_{\mathbf{y}} = \lambda_{PoG} \triangleq \lambda_t + \lambda_{\mathbf{y}'}$, we have our loss function

$$\mathcal{L}_{t-1}^{PoGDiff} = \mathbb{E}_q \left[\frac{\lambda_{\mathbf{y}}}{2} \left\| \mu_{\theta}(\mathbf{x}_t, \mathbf{y}) - \mu_{PoG} \right\|^2 \right] + C.$$
(3)

Here, C is a constant, and μ_{PoG} denotes the mean of the PoG, with the expression defined in Eqn. (2). Then, through derivations based on Gaussian properties, we obtain

$$\mathcal{L}_{t-1}^{PoGDiff} \leq \mathbb{E}_{q} \left[\mathcal{A}(\lambda_{t}) \left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon \right\|^{2} + \mathcal{A}(\lambda_{\mathbf{y}'}) \left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}') \right\|^{2} \right] + C$$
(4)

215 where the function $\mathcal{A}(\lambda) \triangleq \frac{\lambda(1-\alpha_t)^2}{2\alpha_t(1-\bar{\alpha}_t)}$.

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The proof is available in the Appendix. Eqn. (4) in Proposition 3.1 provides a upper bound for the KL divergence (Eqn. (1)) we aim to minimize.

In diffusion model literature (Ho et al., 2020; Rombach et al., 2022), one typically sets $\mathcal{A}(\lambda_t) = 1$ to eliminate the dependency on the time step *t*, and thus Eqn. (4) can be written as¹:

$$\mathcal{L}_{\text{simple}}^{\text{PoGDiff}} = \mathbb{E}_{\mathbf{x}_{0} \sim q(\mathbf{x}_{0}), \epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I}), t \sim \mathcal{U}(1, T)} \left[\left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon \right\|^{2} + \frac{\lambda_{\mathbf{y}'}}{\lambda_{t}} \left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}') \right\|^{2} \right].$$
(5)

For convenience, we rewrite $\frac{\lambda_{\mathbf{y}'}}{\lambda_t} = \frac{\sigma_t^2}{\sigma_{\mathbf{y}'}^2}$. Note that this weight still depends on the time step *t*. Therefore, to be consistent with the DDPM-related literature (Ho et al., 2020; Rombach et al., 2022), we hypothetically define $\sigma_{\mathbf{y}'}^2 = \frac{\sigma_t^2}{\psi[(\mathbf{x}, \mathbf{y}), (\mathbf{x}', \mathbf{y}')]}$ to cancel out the term σ_t^2 , thereby effectively removing the time step dependency; here $\psi[(\mathbf{x}, \mathbf{y}), (\mathbf{x}', \mathbf{y}')]$ denotes the similarity between the two image-text pairs. By shortening the notation $\psi[(\mathbf{x}, \mathbf{y}), (\mathbf{x}', \mathbf{y}')]$ to ψ , we can further rewrite the objective function for PoGDiff as:

$$\mathcal{L}_{\text{simple}}^{\text{PoGDiff}} = \mathbb{E}_{\mathbf{x}_{0} \sim q(\mathbf{x}_{0}), \epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I}), t \sim \mathcal{U}(1, T)} \left[\left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon \right\|^{2} + \psi \left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}') \right\|^{2} \right].$$
(6)

235 3.2.3 Computing the Similarity ψ

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Next, we discuss the choice of ψ in Eqn. (6). Given a image-text dataset \mathcal{D} , the similarities between each image-text pair need to be considered in two parts:

$$\psi \triangleq \psi_{\text{img-sim}} \left(\mathbf{x}, \mathbf{x}' \right) \cdot \psi_{\text{inv-txt-den}} \left(\mathbf{y} \right), \tag{7}$$

where $\psi_{\text{img-sim}}(\mathbf{x}, \mathbf{x}')$ is the similarity between images \mathbf{x} and \mathbf{x}' , and $\psi_{\text{inv-txt-den}}(\mathbf{y})$ is the probability density of the text embedding \mathbf{y} (more details below).

Image Similarity $\psi_{\text{img-sim}}$. For all $\mathbf{x} \sim \mathcal{D}$, we apply a pre-trained image encoder to obtain the latent representations \mathbf{z} . We then calculate the cosine similarities between each \mathbf{z} and select the *k* nearest neighbors with the highest similarity values for all samples in the dataset \mathcal{D} , denoted as $[s_j]_{j=1}^k$, where s_j represents the cosine similarity scores between \mathbf{x} and other images in \mathcal{D} , sorted in descending order. These values are then normalized to produce the weights for each neighbor:

$$v_j = \frac{s_j}{\sum_j s_j},\tag{8}$$

For each data pair (\mathbf{x}, \mathbf{y}) , we then randomly sample a neighboring pair $(\mathbf{x}', \mathbf{y}')$ through from a categorical distribution $Cat([w_j]_{j=1}^k)^2$, i.e., with w_j serving as the probability weight, and compute their image similarity as:

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$$\psi_{\text{img-sim}}\left(\mathbf{x}, \mathbf{x}'\right) \triangleq \max\left(0, s^{a_1 + a_2 \cdot \mathbb{1}\left[\mathcal{I}(\mathbf{x}) \neq \mathcal{I}(\mathbf{x}')\right]}\right),\tag{9}$$

where s denotes the cosine similarity sampled by the weights $\{w_j\}$ defined in Eqn. (8), $\mathbb{1}[\cdot]$ denotes the indicator function, and $\mathcal{I}(\cdot)$ retrieves the class/identity of the current input image; for example, $\mathbb{1}[\mathcal{I}(\mathbf{x}) \neq \mathcal{I}(\mathbf{x}')] = 0$ if \mathbf{x} and \mathbf{x}' are two photos of the same person (e.g., Albert Einstein), and $\mathbb{1}[\mathcal{I}(\mathbf{x}) \neq \mathcal{I}(\mathbf{x}')] = 1$ if \mathbf{x} and \mathbf{x}' are photos of two different persons (e.g., \mathbf{x} is Einstein and \mathbf{x}' Reagan). a_1, a_2 are hyperparameters that control the scale of the similarities.³ The intuition is to compute the image similarity according to both the image content similarity, i.e., s, and identity similarity, i.e., $\mathcal{I}(\mathbf{x})$ and $\mathcal{I}(\mathbf{x}')$.

¹For clarification, our $\mathcal{A}(\lambda_t)$ is equivalent to λ_t in (Ho et al., 2020), with the difference that in our paper, λ refers to the precision of the Gaussian distribution.

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³For example, if the cosine similarity (s) between x and x' is 0.4, and $a_1 = a_2 = 1$: if x and x' are of the same person, the image similarity will be 0.4^1 , whereas if x and x' are not of the same person, the image similarity will be 0.4^2 , which is smaller.



Figure 4: Overview of label distribution for four IT2I datasets. The x-axis corresponds to the identities (i.e., people or individuals).

Inverse Text Densities $\psi_{inv-txt-den}$. Inspired by LDS in DIR (Yang et al., 2021) and the theoretical analysis in VIR (Wang & Wang, 2024), re-weighting the label distribution of an imbalanced dataset can increase the optimization scale for minority classes and reduce the emphasis on majority classes, resulting in better performance under imbalanced conditions. However, both DIR and VIR partition the label space into bins, treating it as a classification problem. This is *not applicable* to our setting because in text-to-image generation, the "label" is actually text embeddings. Instead, we train a variational autoencoder (VAE) on this dataset and then approximate its likelihood $p(\mathbf{y})$ through its evidence lower bound, or ELBO:

$$p(\mathbf{y}) = e^{\log p(\mathbf{y})} \approx e^{\text{ELBO}_{\text{VAE}}(\mathbf{y})}.$$
(10)

The evidence for minority data will be lower than for majority classes. This then motivates our inverse text densities defined as follows:

$$\psi_{\text{inv-txt-den}}\left(\mathbf{y}\right) \triangleq \frac{1}{a_3 \cdot e^{\text{ELBO}_{\text{VAE}}\left(\mathbf{y}\right)}},$$
(11)

where a_3 is a hyperparameter that controls the scale of the inverse text densities. By combining Eqn. (9) and Eqn. (11) to Eqn. (7), we can then compute ψ as follows:

$$\psi = \max\left(0, s^{a_1 + a_2 \cdot \mathbb{1}\left[\mathcal{I}(\mathbf{x}) \neq \mathcal{I}(\mathbf{x}')\right]}\right) \cdot \frac{1}{a_3 \cdot e^{\text{ELBO}_{\text{VAE}}(\mathbf{y})}}$$
(12)

3.2.4 FINAL OBJECTIVE FUNCTION

By collecting all the components discussed above, we arrive at our final training objective:

$$\mathcal{L}_{\text{final}}^{\text{PoGDiff}} = \mathbb{E}_{\mathbf{x}_0 \sim q(\mathbf{x}_0), \epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I}), t \sim \mathcal{U}(1, T)} \left[\left\| \epsilon_{\theta}(\mathbf{x}_t, \mathbf{y}) - \epsilon \right\|^2 + \psi \left\| \epsilon_{\theta}(\mathbf{x}_t, \mathbf{y}) - \epsilon_{\theta}(\mathbf{x}_t, \mathbf{y}') \right\|^2 \right], \quad (13)$$

where ψ is defined in Eqn. (12). Alg. 1 summarizes our algorithm.

4 EXPERIMENTS

4.1 EXPERIMENTAL SETUP

Datasets. To demonstrate the effectiveness of PoGDiff in terms of both accuracy and quality, we evaluate our method on two widely used imbalanced datasets, i.e., AgeDB-IT2I (Moschoglou et al.,

Table 1: P	erformance based on FID score.										
Datasets	AgeDB-IT2I DigiFace-IT2										
Size	Small Medium Large Large										
Metric	$\operatorname{FID}\downarrow$										
Shot	All Few All Few All Few All Few										
VANILLA	14.88 13.72 12.87 12.56 7.67 11.67 7.18 12.23										
CBDM	14.72 14.13 11.63 11.59 7.18 11.12 6.96 12.72										
T2H POGDIFF (OURS)	14.85 13.66 12.79 12.52 7.61 11.64 7.14 12.22 14.15 12.88 10.89 10.64 6.03 10.16 6.84 11.21										

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Table 2: Performance based on D	DINO score.
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Datasets	Ag	DigiFace-IT2I									
Size	Small M	Aedium La	rge Large								
Metric	DINO (cosine similarity) scores ↑										
Shot	All Few A	All Few All	Few All Few								
	0.54 0.09 0. 0.43 0.39 0.	39 0.28 0.34 38 0.11 0.41 42 0.29 0.37 69 0.56 0.66	0.26 0.34 0.16 0.26 0.44 0.36								

334 Table 3: Performance on AgeDB-IT2I based on Table 4: Performance on AgeDB-IT2I based 335 **human evaluation.** The evaluation is a binary de- on GPT-40 evaluation. The scores are from 0 to 336 cision: the image is either judged as representing 10, with higher scores indicating the individual the same individual (score 1.0) or not (score 0.0). resembles the well-known person. 337

Size	Small Medium Large	Size Small Medium Large
Metric	Human Score ↑	Metric GPT-40 Evaluation ↑
Shot	All Few All Few All Few	Shot All Few All Few All Fe
VANILLA	0.50 0.00 0.66 0.32 0.60 0.20	VANILLA 5.20 3.20 4.30 2.90 4.90 3.
CBDM	0.50 0.00 0.44 0.08 0.56 0.12	CBDM 4.50 1.10 1.30 1.00 3.10 1.
T2H	0.50 0.00 0.66 0.32 0.60 0.20	T2H 5.50 3.10 4.60 3.00 4.70 3.
POGDIFF (OURS)	1.00 1.00 0.96 0.92 0.84 0.68	POGDIFF (OURS) 9.10 8.40 8.80 8.20 8.50 8.

347 2017) and DigiFace-IT2I (Bae et al., 2023). Note that our method is designed for fine-tuning. Therefore our setup does not require large-scale, long-tailed human datasets. Instead, we sample 348 from these datasets, as long as they meet the following criteria: (1) the dataset must be long-tailed, (2) 349 traditional methods must fail to recognize the minority classes, and (3) there must be a distinguishable 350 difference between the majority and minority classes (e.g., we prefer visual distinctions between the 351 two groups to better highlight the impact of our method). Fig. 4 shows the label density distribution 352 of these datasets, and their level of imbalance⁴. 353

AgeDB-IT21: AgeDB-IT2I is constructed from the AgeDB dataset (Moschoglou et al., 2017). For 354 each image x in AgeDB, we passed it through the pretrained LLaVA-1.6-7b model (Liu et al., 2024) 355 to generate textual captions \tilde{y} . Since the identities in AgeDB are well-known individuals that the 356 pretrained SDv1.5 (Rombach et al., 2022) might have encountered during pre-training, we masked 357 the true names and replaced them with generic, random names, leading to a new caption y. For 358 example, we replace "Albert Einstein" in the caption with a random name "Lukas". Finally, we 359 collect all (\mathbf{y}, \mathbf{x}) pairs to form our AgeDB-IT2I dataset. 360

Additionally, given that the identities (i.e., people or individuals) in AgeDB are well-known figures, 361 we sampled from AgeDB to create three datasets for comprehensive analysis: AgeDB-IT2I-L (large), 362 AgeDB-IT2I-M (medium), and AgeDB-IT2I-S (small). Specifically: 363

- AgeDB-IT21-L (large). This dataset consists of 976 images across 223 identities, with each majority class containing 30 images and each minority class containing 2 images.
- AgeDB-IT2I-M (medium). This dataset consists of 100 images across 10 identities, with each majority class containing 30 images and each minority class containing 2 images.
- AgeDB-IT2I-S (small). This dataset contains 32 images across 2 identities, where each majority class consists of 30 images and each minority class consists of 2 images.

DigiFace-IT21-L: DigiFace-IT21-L is derived from the DigiFace dataset (Bae et al., 2023). It contains 372 985 images across 179 identities, where each majority class consists of 30 images and each minority 373 class consists of 2 images. We use a process similar to AgeDB-IT2I to collect text-image pairs, forming this DigiFace-IT2I dataset. 374

375 Baselines. We employ Stable Diffusion v1.5 (Rombach et al., 2022) as the backbone diffusion model. 376 As this is the first work to explore imbalanced text-to-image (IT2I) diffusion models with **natural**

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⁴Our datasets actually contain sparse datasets; more details can be found in Appendix C.4.



Implementation Details. For both baselines and our model, we used the same hyper-parameter settings, specifically

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• AgeDB-IT2I-L & DigiFace-IT2I-L. The learning rate was set to 1×10^{-5} , with a maximum of 12,000 training steps. The effective batch size per GPU was 32, calculated as 8 (Batch Size) \times 4 (Gradient Accumulation Steps).

• AgeDB-IT2I-M & AgeDB-IT2I-S. The learning rate was set to 1×10^{-5} , with a maximum of 6,000 training steps. The effective batch size per GPU was 8, calculated as 8 (Batch Size) \times 1 (Gradient Accumulation Step).

4.2 Results

444 We report the performance of different methods in terms of FID score, human evaluation score, 445 GPT-40 score, and DINO score in Table 1, Table 2, Table 3 and Table 4, respectively⁵. Across 446 all tables, we observe that our PoGDiff consistently outperforms all baselines. Notably, PoGDiff 447 demonstrates significant improvements, especially in few-shot scenarios (i.e., for minority classes). 448 It is also worth noting that CBDM (Qin et al., 2023) performs extremely poorly on AgeDB-IT2I-S 449 and AgeDB-IT2I-M datasets. This is because their method samples text conditions from the entire space, which may work in one-hot class settings, but in our context (natural text conditions), this 450 sampling technique misguides the model during training. In addition, for each method, we report the 451 performance on low-density classes in AgeDB-IT2I-L in Fig. 5. Across each column, the individual 452 names are Albert Einstein, JW Marriott, J.P. Morgan, Edward G. Robinson, Larry Ellison, and Luise 453 Rainer, respectively. The results show that our PoGDiff achieves significantly better accuracy and 454 quality for tail classes. 455

Note that one of our primary objectives is to generate accurate images of the same individual while
ensuring facial consistency. Therefore diversity can sometimes be harmful. For example, given a
text input of "Einstein", generated images with high diversity would generate both male and females
images; this is obviously incorrect. Therefore it is important to strike a balance between diversity
and accuracy, a goal that our PoGDiff achieves.

461 Specifically, as shown in Fig. 5:

- First Three Columns of SDv1.5, CBDM, and PoGDiff: In these cases, the training dataset contains only two images per person⁶. With such limited data, it is impossible to introduce meaningful diversity.
 - SDv1.5 fails to generate accurate images altogether in this scenario.
 - While CBDM might appear to produce the "diversity" you mentioned, it does so incorrectly, as it generates an image of a woman when the target is Einstein.
 - In contrast, our PoGDiff can successfully generate accurate images (e.g., Einstein images in Column 1) while still enjoying sufficient diversity.
 - Fourth and Fifth Columns: Here, the training dataset contains a medium number of images per person (5–7 images). Under these conditions:
 - SDv1.5 can generate accurate representations of individuals, but its outputs lack diversity.
 - CBDM, on the other hand, introduces "diversity" but consistently generates incorrect results.
 - In contrast, our method produces accurate images of the target individual while demonstrating greater diversity than SDv1.5.
 - Sixth Column: In this case, the training dataset includes 30 images per person.
 - SDv1.5 generates accurate images but with nearly identical expressions, i.e., poor diversity.
 - CBDM still fails to generate accurate depictions of the individual.
 - In contrast, PoGDiff successfully generates accurate images while maintaining diversity.

In summary, typical diversity evaluation in diffusion model evaluations, such as generating multiple
types of trees for a "tree" prompt, is not the focus of our setting and may even be misleading. In
our setting, the key is to balance accuracy and diversity.

⁵Note that the CLIP score is not applicable in our setting. Specifically, our text prompts are predominantly
human names. However, CLIP is primarily trained on common objects, not human names; therefore the CLIP
score can not be use to compute matching scores between images and human names.

⁶More details are included in the Appendix.

486 5 ABLATION STUDY

To verify the effectiveness of each component in the second term in our PoGDiff final objective function from Eqn. (12), we report the accuracy of our proposed PoGDiff after removing the y' (i.e., same as Vanilla model), the Image Similarity term $\psi_{img-sim}$, and/or the Inverse Text Densities term $\psi_{inv-txt-den}$ in Table 5 for AgeDB-IT2I-L. The results show that removing either term leads to a performance drop, confirming the importance of both terms in our PoGDiff.

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6 LIMITATIONS

496 Datasets. Our method relies heavily497 on "borrowing" the statistical strength

498 of neighboring samples from minor-499 ity classes, making the results sensi-500 tive to the size of the minority class. (i.e., in our assumption we require at 501 least 2 for each minority class). In ad-502 dition, while our AgeDB-IT2I-small 503 and AgeDB-IT2I-medium are actu-504 ally the sparse dataset, the cardinality 505 remains limited in our experiments. 506 Therefore, how to address IT2I prob-

Table 5: Ablation Studies.

Datasets			Age	DB-I7	72 I -I 2	rge		
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Size	FI	D↓	Hum	an ↑	GPT	-4o †	DIN	iO ↑
Shot	All	Few	All	Few	All	Few	All	Few
w/o \mathbf{y}' (Vanilla)	7.67	11.67	0.60	0.20	4.90	3.60	0.34	0.25
w/o $\psi_{\text{IMG-SIM}}$	6.41	10.49	0.84	0.68	8.40	7.60	0.57	0.46
w/o $\psi_{\text{inv-txt-den}}$	6.35	10.43	0.84	0.68	8.20	7.80	0.64	0.51
POGDIFF (OURS)	6.03	10.16	0.84	0.68	8.50	8.00	0.66	0.52

lem under this settings are interesting directions.

Models. Our method is a general fine-tuning approach designed for datasets that the Stable Diffusion 509 (SD) model has not encountered during pre-training. As shown in Fig. 1, color deviation is very 510 common and is a known issue when one fine-tunes diffusion models (as also mentioned in (Song 511 et al., 2020b)); for example, we can observe similar color deviation in both baselines (e.g., CBDM 512 and Stable Diffusion v1.5) and our PoGDiff. This can be mitigated using the exponential moving 513 average (EMA) technique (Song et al., 2020b); however, this is orthogonal to our method and is 514 outside the scope of our paper. Moreover, as shown in Fig. 5, the baseline Stable Diffusion also 515 suffers from this issue. Besides, exploring PoGDiff's performance when training from scratch is also 516 an interesting direction for future work.

Evaluation. Our goal is to adapt the pretrained diffusion model to a specific dataset; therefore the evaluation should focus on the target dataset rather than the original dataset used during pre-training. For example, when a user fine-tunes a model on a dataset of employee faces, s/he is not interested in how well the fine-tuned model can generate images of "tables" and "chairs". Evaluating the model's performance on the original dataset used during pre-training would be an intriguing direction for future work, but it is orthogonal to our proposed PoGDiff and out of the scope of our paper.

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7 CONCLUSIONS

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In this paper, we propose a general fine-tuning approach called PoGDiff to address the performance drop that occurs when fine-tuning on imbalanced datasets. Instead of directly minimizing the KL divergence between the predicted and ground-truth distributions, PoGDiff replaces the ground-truth distribution with a Product of Gaussians (PoG), constructed by combining the original groundtruth targets with the predicted distribution conditioned on a neighboring text embedding. Looking ahead, an interesting avenue for future research would be to explore more innovative techniques for re-weighting minority classes (as discussed in Sec. 6), particularly within the constraints of: (1) long-tailed generation settings, as opposed to recognition tasks, and (2) natural text prompts rather than one-hot class labels.

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A PROOFS FOR PROPOSITION 3.1

Proposition A.1. Assume $\lambda_{\mathbf{y}} = \lambda_{PoG} \triangleq \lambda_t + \lambda_{\mathbf{y}'}$, we have our loss function

$$\mathcal{L}_{t-1}^{PoGDiff} = \mathbb{E}_q \left[\frac{\lambda_{\mathbf{y}}}{2} \left\| \mu_{\theta}(\mathbf{x}_t, \mathbf{y}) - \mu_{PoG} \right\|^2 \right] + C.$$
(14)

Here, C is a constant, and μ_{PoG} denotes the mean of the PoG, with the expression defined in Eqn. (2). Then, through derivations based on Gaussian properties, we obtain

$$\mathcal{L}_{t-1}^{PoGDiff} \leq \mathbb{E}_{q} \left[\mathcal{A}(\lambda_{t}) \left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon \right\|^{2} + \mathcal{A}(\lambda_{\mathbf{y}'}) \left\| \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}) - \epsilon_{\theta}(\mathbf{x}_{t}, \mathbf{y}') \right\|^{2} \right] + C$$
(15)

715 where the function $\mathcal{A}(\lambda) \triangleq \frac{\lambda(1-\alpha_t)^2}{2\alpha_t(1-\bar{\alpha}_t)}$.

Proof. To prove the above inequality, we need to prove the following lemma.

Lemma A.1. Assume $\lambda_{\mathbf{y}} = \lambda_{PoG} \triangleq \lambda_t + \lambda_{\mathbf{y}'}$, and for simplicity we shorten the notation from $\epsilon_{\theta}(\mathbf{x}_t, \mathbf{y})$ and $\mu_{\theta}(\mathbf{x}_t, \mathbf{y})$ to $\epsilon_{\theta}(\mathbf{y})$ and $\mu_{\theta}(\mathbf{y})$, respectively. Then we have

$$\frac{1}{2}\lambda_{t}\left(\mu_{\theta}\left(\mathbf{y}\right)-\mu_{t}\right)^{2}+\frac{1}{2}\lambda_{\mathbf{y}'}\left(\mu_{\theta}\left(\mathbf{y}\right)-\mu_{\theta}\left(\mathbf{y}'\right)\right)^{2}\geq\frac{1}{2}\lambda_{\mathbf{y}}\left(\mu_{\theta}\left(\mathbf{y}\right)-\mu_{PoG}\right)^{2}$$
(16)

Proof. By the definition of Gaussian property, we have

$$\begin{split} &\frac{1}{2}\lambda_{t}\left(\mu_{\theta}(\mathbf{y})-\mu_{t}\right)^{2}+\frac{1}{2}\lambda_{\mathbf{y}'}\left(\mu_{\theta}(\mathbf{y})-\mu_{\theta}(\mathbf{y}')\right)^{2} \\ &=\frac{[\mu_{\theta}(\mathbf{y})]^{2}-2\mu_{t}\mu_{\theta}(\mathbf{y})+\mu_{t}^{2}}{2\lambda_{t}^{-1}}+\frac{[\mu_{\theta}(\mathbf{y}')]^{2}-2\mu_{\theta}(\mathbf{y}')\mu_{\theta}(\mathbf{y})+[\mu_{\theta}(\mathbf{y}')]^{2}}{2\lambda_{\mathbf{y}'}^{-1}} \\ &=\frac{\left(\lambda_{t}^{-1}+\lambda_{\mathbf{y}'}^{-1}\right)[\mu_{\theta}(\mathbf{y})]^{2}-2\left(\frac{\mu_{t}}{\lambda_{\mathbf{y}'}}+\frac{\mu_{\theta}(\mathbf{y}')}{\lambda_{t}}\right)\mu_{\theta}(\mathbf{y})+\frac{\mu_{t}^{2}}{\lambda_{\mathbf{y}'}}+\frac{[\mu_{\theta}(\mathbf{y}')]^{2}\lambda_{\mathbf{y}'}}{\lambda_{t}} \\ &=\frac{[\mu_{\theta}(\mathbf{y})]^{2}-2\left(\frac{\mu_{t}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]\lambda_{\mathbf{y}'}}{\lambda_{t}+\lambda_{\mathbf{y}'}}\right)\mu_{\theta}(\mathbf{y})+\frac{\mu_{t}^{2}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]^{2}\lambda_{\mathbf{y}'}}{\lambda_{t}+\lambda_{\mathbf{y}'}} \\ &+\frac{\left[\frac{\mu_{t}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]\lambda_{\mathbf{y}'}}{\lambda_{t}+\lambda_{\mathbf{y}'}}\right]^{2}}{\frac{2}{\lambda_{t}}+\lambda_{\mathbf{y}'}} \\ &=\frac{\left(\mu_{\theta}(\mathbf{y})-\frac{\mu_{t}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]\lambda_{\mathbf{y}'}}{\lambda_{t}+\lambda_{\mathbf{y}'}}\right)^{2}}{\frac{2}{\lambda_{t}}+\lambda_{\mathbf{y}'}} +\frac{\left(\mu_{t}^{2}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]^{2}\lambda_{\mathbf{y}'})(\lambda_{t}+\lambda_{\mathbf{y}'})-(\mu_{t}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]\lambda_{\mathbf{y}'})^{2}}{2(\lambda_{t}+\lambda_{\mathbf{y}'})} \\ &=\frac{\left(\mu_{\theta}(\mathbf{y})-\frac{\mu_{t}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]\lambda_{\mathbf{y}'}}{\lambda_{t}+\lambda_{\mathbf{y}'}}\right)^{2}}{\frac{2}{\lambda_{t}}+\lambda_{\mathbf{y}'}}} +\frac{\left(\mu_{t}^{2}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]^{2}\lambda_{\mathbf{y}'})(\lambda_{t}+\lambda_{\mathbf{y}'})-(\mu_{t}\lambda_{t}+[\mu_{\theta}(\mathbf{y}')]\lambda_{\mathbf{y}'})^{2}}{2(\lambda_{t}+\lambda_{\mathbf{y}'})} \\ &=\frac{1}{2}\lambda_{\mathbf{y}}\left(\mu_{\theta}(\mathbf{y})-\mu_{\mathrm{PoG}}\right)^{2} +\frac{\lambda_{t}\lambda_{\mathbf{y}'}(\mu_{t}-\mu_{\theta}(\mathbf{y}'))^{2}}{2(\lambda_{t}+\lambda_{\mathbf{y}'})} \\ &\geq\frac{1}{2}\lambda_{\mathbf{y}}\left(\mu_{\theta}(\mathbf{y})-\mu_{\mathrm{PoG}}\right)^{2}. \end{split}$$

Thus we complete the proof.

 $\frac{1}{2}\lambda_{\mathbf{y}} \left\| \mu_{\theta}(\mathbf{y}) - \mu_{\text{PoG}} \right\|^2 \equiv \frac{1}{2}\lambda_{\mathbf{y}} \left(\mu_{\theta}(\mathbf{y}) - \mu_{\text{PoG}} \right)^2$

From Lemma A.1, we can derive

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where the function $\mathcal{A}(\lambda) \triangleq \frac{\lambda(1-\alpha_t)^2}{2\alpha_t(1-\bar{\alpha}_t)}$, and the last equivalence is because the transform from $\mu_{\theta}(\cdot)$ to $\epsilon_{\theta}(\cdot)$.

 $\leq \quad \frac{1}{2}\lambda_t \left(\mu_{\theta}(\mathbf{y}) - \mu_t\right)^2 + \frac{1}{2}\lambda_{\mathbf{y}'} \left(\mu_{\theta}(\mathbf{y}) - \mu_{\theta}(\mathbf{y}')\right)^2$

 $\equiv \frac{1}{2}\lambda_t \|\mu_{\theta}(\mathbf{y}) - \mu_t\|^2 + \frac{1}{2}\lambda_{\mathbf{y}'} \|\mu_{\theta}(\mathbf{y}) - \mu_{\theta}(\mathbf{y}')\|^2$

 $\equiv \mathcal{A}(\lambda_t) \left\| \epsilon_{\theta}(\mathbf{y}) - \epsilon \right\|^2 + \mathcal{A}(\lambda_{\mathbf{y}'}) \left\| \epsilon_{\theta}(\mathbf{y}) - \epsilon_{\theta}(\mathbf{y}') \right\|^2,$

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B DETAILS FOR EVALUATION

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772 In this section, we provide details on our evaluation procedures.

FID Score. For each identity, we collect all images from the original AgeDB or DigiFace datasets
as the *true image set*. Then, In *all-shot* evaluation, for AgeDB-IT2I-S and AgeDB-IT2I-M, we
generate 100 images per identity as the *fake image set*, and for AgeDB-IT2I-L and DigiFace-IT2I-L,
we generate 20 images per identity as the *fake image set*. In *few-shot* evaluation, we we generate
500 images per identity as the *fake image set*. For all generations, we employ the DDIM sampling
technique (Song et al., 2020a) with 50 steps. The prompt used during generation is "An image of
{p}." where "p" is the name of the identity (e.g., Albert Einstein).

Human & GPT-40 Feedback. For each minority identity, we generate 5 images using DDIM 781 sampling (Song et al., 2020a) with 50 steps. We then ask 10 people to evaluate whether the images depict the same person (scored as 1.0) or not (scored as 0.0). Additionally, for each image, we ask 782 the GPT-40 model to rate the similarity on a scale from 1 to 10. The prompt used during generation 783 is "An image of $\{p\}$." where "p" is the name of the identity. The text prompt using for GPT-40 784 model is "It is mandatory to give a score that how close the person in the image to a well-known 785 individual. A score of 10.0 means they are exactly the same person, while a score of 0.0 means they 786 are definitely not the same person. How close you think the person in the image is to 'p-true'." where 787 "p-true" denotes the real name (well-known name) in AgeDB. Note that the GPT-40 model might 788 occasionally refuse to provide a score, and you may need to repeat and compel it to give a rating. For 789 each image, we collect 10 scores from the GPT-40 model and report the average rating. 790

Evaluating Image Similarities. We collect samples that are outside our training dataset (e.g., AgeDB-T2I-L) but belong to the original dataset (e.g., AgeDB). Using the same prompt, we generate the corresponding images. A pre-trained DINOv2 model (Caron et al., 2021) is then applied to extract latent features, and cosine similarities are calculated.

- C DISCUSSION
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C.1 PROBLEM SETTINGS

We would like to clarify that our paper focuses on a setting different from works like Dream-Booth (Ruiz et al., 2023), and our focus is not on diversity, but on finetuning a diffusion model on an imbalanced dataset. Specifically:

- Different Setting from Custom Techniques like DreamBooth (Ruiz et al., 2023), CustomDiffusion (Kumari et al., 2023) and PhotoMaker (Li et al., 2024). Previous works like CustomDiffusion and PhotoMaker focus on adjusting the model to generate images with a single object, e.g., a specific dog. In contrast, our PoGDiff focuses finetuning the diffusion model on an entire data with many different objects/persons simultaneously. They are very different settings and are complementary to each other.
- **Diversity.** Note that while our PoG can naturally generate images with diversity, diversity is actually **not** our focus. Our goal is to fine-tune a diffusion model on an imbalanced dataset.



Figure 6: Example generated images from different methods. Our PoGDiff outperforms the baselines in both generation accuracy and quality. Regarding the ground truth (GT), the training set for the minority class (left two columns) contains only 2 images per individual, whereas the majority class has more than 10 samples per individual.

For example, PoGDiff can fine-tune a diffusion model on an imbalanced dataset of employee faces so that the diffusion model can generate new images that match each employee's identity. In this case, we are more interested in "faithfulness" rather than "diversity".

C.2 UNDERSTANDING FIG. 5

To further emphasize our results, note that one of our primary objectives is to generate accurate images of the same individual while ensuring facial consistency. Therefore **diversity can be harmful**. For example, given a text input of "Einstein", generated images with high diversity would generate both male and females images; **this is obviously incorrect**. Therefore it is important to strike a balance between **diversity** and **accuracy**, a goal that our PoGDiff achieves.

Fig. 6 (which contains images from Column 1, 2, and 6 for each method in Fig. 5) provides a clearer comparison with the training images. Specifically:

- Ground-Truth (GT) Images: We show the ground-truth images on the right-most 3 columns.
- Column 1 and 2 of SDv1.5, CBDM, PoGDiff, and GT: In these cases, the training dataset contains only two images per person. With such limited data, it is impossible to introduce meaningful diversity.
 - SDv1.5 fails to generate accurate images altogether in this scenario.
- While CBDM might appear to produce the "diversity" you mentioned, it does so incorrectly, as it generates an image of a woman when the target is Einstein (we circled those wrong samples in first column in Fig. 6).



• Case 3: y' is very close to y. When y' is close to y, the reweighting can be approximated as: $\alpha y + (1 - \alpha)y' \approx y + (1 - \alpha)(y' - y)$. Since y' is nearly identical to y, this 

Figure 9: TSNE visualization for all the methods for an example individual in the Age-DB-IT2Mlarge dataset.

effectively introduces a small weighted noise term $(1-\alpha)(\mathbf{y}'-\mathbf{y})$ into y. In our preliminary experiments, this additional noise degraded the performance compared to the original baseline results.

Based on these observations, direct smoothing of text embeddings appears ineffective and may even harm performance in some cases.

OUR DATASET COVERS DIFFERENT LEVELS OF SPARSITY C.4

Our AgeDB-IT2M-small and AgeDB-IT2M-medium datasets are actually very sparse and are meant for evaluate the sparse data. For example, the AgeDB-IT2M-small only contains images from 2 persons, it is therefore a very sparse data setting, compared to AgeDB-IT2M-large with images across 223 persons. Fig. 8 shows the bar plot version for our datasets, while sparse settings are not our primary focus, we agree that addressing imbalanced image generation in such setting is an interesting and valuable direction, and we have included a discussion about this in the limitations section of the paper.

C.5 DISCUSSION ON FID

It is important to note that the FID score measures only the distance between Gaussian distributions of ground-truth and generated images, relying solely on mean and variance. As a result, it does not fully capture the nuances of our task. This is why we include additional evaluation metrics such as DINO Score, Human Score, and GPT-40 Score, to comprehensively verify our method's superiority (as shown in Table 2, Table 3 and Table 4).

Additional Experiments: Limitation of FID. In addition, we have added a figure showcasing a t-SNE visualization for a minority class as an example, as shown in Fig. 9, to further illustrate the limitation of FID we mentioned above. As shown in the figure:

- There are two ground-truth IDs (i.e., two ground-truth individuals) in the training set.
- Our PoGDiff can successfully generate images similar to these two ground-truth ID while maintaining diversity.
- • All baselines, including CBDM, fail to generate accurate images according to the ground-truth IDs. In fact most generated images from the baselines are similar to other IDs, i.e., generating the facial images of wrong individuals.
 - These results show that:

- Our PoGDiff significantly outperforms the baselines.
 - FID fails to capture such improvements because it depends only on the mean and variance of the distribution, losing a lot of information during evaluation.

D MORE BASELINES

In this section, we also include one more work related to our baseline CBDM (Qin et al., 2023); they are both equivalent to direct reweighting/resampling. We did not include (Zhang et al., 2024) as a baseline because it is not directly applicable to our setting. Specifically, (Zhang et al., 2024) relies on the class frequency, which is not available in our setting. Therefore, we adapted this method to our settings by using the density for each text prompt embedding to serve as the class frequency in (Zhang et al., 2024). Results shown in Fig. 7 show that it performs even worse than CBDM, and it performs similar to directly fine-tuning a SD model.

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E MORE EVALUATION METRIC: RECALL SCORE.

FID Measures Both ID Consistency

and Diversity. We could like to clarify that our Fréchet Inception Distance
(FID) is computed for each ID separately, and the final FID score in the tables (e.g., Table 1) is the average FID
over all IDs. Therefore FID measures
both ID consistency and diversity.

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997 To see why, note that the FID score measures the distance between two Gaussian distributions, where the *mean* of the Gaussian represents the *identity (ID)* and the *variance* represents the *diversity*. For example, the

 Table 6: Recall for the AgeDB-IT2I Dataset. See the detailed definition of recall in Appendix E.

Size	Sn	nall	Med	lium	La	rge					
Metric	Recall Score ↑										
Shot	All	Few	All	Few	All	Few					
VANILLA		0.000									
CBDM	0.267	0.000	0.159	0.083	0.138	0.100					
T2H	0.017	0.000	0.104	0.167	0.196	0.200					
POGDIFF (OURS)	0.800	1.000	0.517	0.642	0.435	0.540					

mean of the ground-truth distribution represents the embedding position of the ground-truth ID,
 while the *variance* of the ground-truth distribution represents the *diversity* of ground-truth images.
 Similarly, the *mean* of the generated-image distribution represents the embedding position of the
 generated-image ID, while the *variance* of the generated-image distribution represents the *diversity* of generated images. A lower FID score indicates that the generated-image distribution more closely
 matches the ground truth distribution in terms of both ID and diversity.

1008 Results Related to Diversity. Currently,

- **PoGDiff's Superior FID Performance.** as shown in Table 1, we demonstrate that PoGDiff achieves a lower FID score, particularly in few-shot regions (i.e., minorities). This suggests that the images generated by our method capture a broader range of variations present in the training dataset, such as **backgrounds or facial angles**.
 - **PoGDiff's Visualization.** As shown in Fig. 6:
 - For Einstein (Column 1 for each method), the training dataset includes two face angles and two hairstyles. Our generated results successfully cover these attributes.
 - For JW Marriott (Column 2 for each method), the training dataset has only one face angle. Correspondingly our results focus on generating subtle variations in facial expressions with only one angle, as expected.
 - For the majority group (Column 3 for each method), our results clearly show that the generated images cover a wider range of diversity while maintaining ID consistency.

Additional Experiments on Recall (a New Metric). To better evaluate the superiority of our PoGDiff, we propose a new metric, "recall".

• **Recall in the Context of Image Generation: "Correct Image" and "Covered Image".** For each generated image, we classify it as a "correct image" if its distance to at least one

ground-truth (GT) image is below a predefined threshold. For instance, suppose we have two
training-set images for Einstein, denoted as x_1 and x_2 . A generated image x_q is a "correct
images if the cosine similarity between x_g and either x_1 or x_2 is above some threshold (e.g.,
we set to 0.9 here). For example, if the cosine similarity x_q and x_1 is larger than 0.9, we
say that x_g is a "correct image", and that x_1 is a "covered image". Intuitively, a training-set
image (e.g., x_1) is covered if a diffusion model is capable of generating a similar image.
• Formal Definition for Recall. Formally, for each model, we compute the Recall per ID as
follows:
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$\operatorname{Recall} = \frac{1}{c} \sum_{i=1}^{c} \frac{\operatorname{number of unique covered images for ID i}}{\operatorname{number of images for ID i in the training set}} $ (17)
$c \underset{i=1}{\frown}$ number of images for ID i in the training set
where c is the number of IDs in a training set.
-
• Cosine Similarity between Images. Note that in practice, we compute the cosine similarity between DINO embeddings of images rather than raw pixels.
• Analysis. This metric evaluates the generational diversity of a model. For example, if the training detect contains two distinct images of Einstein <i>m</i> and <i>m</i> and <i>n</i> model generates
training dataset contains two distinct images of Einstein, x_1 and x_2 , and a model generates only images resembling x_1 , the recall in this case would be 0.5. While the model may
achieve high accuracy in terms of facial identity (Table 3 and Table 4), it falls short in
diversity because it fails to generate images resembling x_2 . In contrast, if a model generates
images that cover both x_1 and x_2 the recall for this ID will be 1; for instance, if the model
generates 10 images for Einstein, where 6 of them resemble x_1 and 4 of them resemble x_2 ,
the recall would be 1, indicating high diversity and coverage.
Additional Results in Terms of Recall. Table 6 shows the recall for different methods on three
datasets, AgeDB-IT2I-small, AgeDB-IT2I-medium, and AgeDB-IT2I-large. These results show that
our PoGDiff achieves much higher recall compared to all baselines, demonstrating its impressive
diversity.
Additional Details for AgeDB-IT2I-small in Table 6. For AgeDB-IT2I-small, there are two IDs, one "majority" ID with 30 images and one minority ID with 2 images.
• For VANILLA and T2H , the recall for the majority ID and the minority ID is $1/30$ and $0/2$,
• For VARIELA and 12H, the recall for the majority 1D and the minority 1D is 1/30 and 0/2, respectively. Therefore, the average recall score is $0.5 * 1/30 + 0.5 * 0/2 \approx 0.0167$.
• For CBDM , the recall for the majority ID and the minority ID is $16/30$ and $0/2$, respectively. Therefore, the average recall score is $0.5 * 16/30 + 0.5 * 0/2 \approx 0.2667$.
• For PoGDiff (Ours) , the recall for the majority ID and the minority ID is $18/30$ and $2/2$,
respectively. Therefore, the average recall score is $0.5 * 18/30 + 0.5 * 2/2 = 0.8$.
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F MORE DATASETS
We have included an additional dataset, VGGFace, for evaluation. Specifically, we constructed a

We have included an additional dataset, VGGFace, for evaluation. Specifically, we constructed a subset from VGGFace2 (Cao et al., 2018), named VGGFace-IT2I-small. This is a sparse dataset consisting of two individuals: the majority group contains 30 images, while the minority group contains only 2 images.

1070 The results shown in Table 7, Table 8, Table 9, Table 10 and Table 11, below demonstrate that our
1071 PoGDiff consistently outperform all baselines, highlighting its robustness and superior performance
1072 even on imbalanced and sparse datasets.

Table 7: Performance based on FID score.											
Datasets		1	AgeDE	B-IT2I		DigiFace-IT2I VGGFace-IT2I					
Size	Sma	all	Med	lium	rge	L	arge	Small			
Metric	FID↓										
Shot	All	Few	All	Few	All	Few	All	Few	All	Few	
VANILLA	14.88	13.72	12.87	12.56	7.67	11.67	7.18	12.23	14.18	12.73	
CBDM	14.72	14.13	11.63	11.59	7.18	11.12	6.96	12.72	13.85	13.21	
T2H	14.85	13.66	12.79	12.52	7.61	11.64	7.14	12.22	14.16	12.74	
POGDIFF (OURS)	14.15	12.88	10.89	10.64	6.03	10.16	6.84	11.21	13.68	11.11	

Table 7. Performance based on FID score

Table 8: Performance based on DINO score.

Datasets	AgeDB-IT2I						DigiFace-IT2I VGGFace-IT2I					
Size	Small Medium Large						L	arge	S	Small		
Metric	DINO (cosine similarity) scores ↑											
Shot	All	Few	All	Few	All	Few	All	Few	All	Few		
VANILLA	0.42	0.37	0.39	0.28	0.34	0.25	0.42	0.36	0.49	0.36		
CBDM	0.54	0.09	0.38	0.11	0.41	0.26	0.34	0.16	0.52	0.06		
T2H	0.43	0.39	0.42	0.29	0.37	0.26	0.44	0.36	0.48	0.37		
POGDIFF (OURS)	0.77	0.73	0.69	0.56	0.66	0.52	0.64	0.49	0.84	0.79		

Table 9: Performance on AgeDB-IT2I based on Table 10: Performance on AgeDB-IT2I based human evaluation. The evaluation is a binary de- on GPT-40 evaluation. The scores are from 0 to cision: the image is either judged as representing 10, with higher scores indicating the individual

resembles the well-known person

the same many	resembles the wen-known person.												
Datasets	4	AgeDB-IT2	I	VGC	Face-IT2I	Datasets		Ag	eDB-IT	2I		VGG	Face-IT2I
Size	Small	Medium	Large		Small	Size	Sm	all M	Medium	La	rge	5	Small
Metric	Metric Human Score ↑						Metric GPT-40 Evaluation ↑						
Shot	All Few	All Few	All Fe	w All	Few	Shot	All	Few A	All Few	All	Few	All	Few
VANILLA	0.50 0.00	0.66 0.32	0.60 0.2	0 0.50	0.00	VANILLA	5.20	3.20 4	.30 2.90	4.90	3.60	6.00	3.60
CBDM	0.50 0.00	0.44 0.08	0.56 0.1	2 0.50	0.00	CBDM	4.50	1.10 1	.30 1.00	3.10	1.70	4.67	1.33
T2H	0.50 0.00	0.66 0.32	0.60 0.2	0 0.50	0.00	T2H	5.50	3.10 4	.60 3.00	4.70	3.90	6.05	3.80
POGDIFF (OURS)	1.00 1.00	0.96 0.92	0.84 0.6	8 1.00	1.00	POGDIFF (OURS)	9.10	8.40 8	.80 8.20	8.50	8.00	7.90	9.60

1113	the same individual (score 1.0) or not (score 0.0)				
1114	Datasets	AgeDB-IT2I VGGFace-IT2I			
1115	Size	Small Medium Large Small			
1116	Metric	Human Score ↑			
1117	Shot	All Few All Few All Few All Few			
1118	VANILLA	0.50 0.00 0.66 0.32 0.60 0.20 0.50 0.00			
1119	CBDM	0.50 0.00 0.44 0.08 0.56 0.12 0.50 0.00			
1120	T2H	0.50 0.00 0.66 0.32 0.60 0.20 0.50 0.00			
1120	POGDIEF (OURS)				

Table 11: Recall for the AgeDB-IT2I Dataset. See the detailed definition of recall in Appendix E.

Datasets	Age	DB-IT2I	VGGFace-IT2I	
Size	Small M	Iedium Large	Small	
Metric	Recall Score ↑			
Shot	All Few Al	1 Few All Fe	ew All Few	
VANILLA	0.017 0.000 0.10	04 0.167 0.196 0.2	200 0.033 0.000	
CBDM	0.267 0.000 0.15	59 0.083 0.138 0.1	00 0.233 0.000	
T2H	0.017 0.000 0.10	04 0.167 0.196 0.2	200 0.033 0.000	
POGDIFF (OURS)	0.800 1.000 0.5	17 0.642 0.435 0.5	540 0.767 1.000	