Ultrafast Source Mask Optimization via Conditional Discrete Diffusion
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Abstract—Source mask optimization (SMO) is vital for mitigating lithography imaging distortions caused by shrinking critical dimensions in integrated circuit fabrication. However, the computational intensity of SMO, involving multiple integrals in Abbe’s theory, hinders its widespread adoption and advancement. In this article, we present Diff-SMO, a highly efficient and accurate SMO framework with a primary emphasis on enhancing source optimization techniques. Previous research was confined to mask optimization acceleration due to the constraints of the academia lithography model. Diff-SMO extends the scope of optimization by concurrently refining the intricate interplay between the source and mask. We first develop a GPU-accelerated lithography simulator grounded in Abbe’s theory, enabling full GPU acceleration throughout the SMO process. Furthermore, we propose a discrete diffusion model for generating quasi-optimal sources, significantly improving computational efficiency. Our experimental results demonstrate exceptional imaging fidelity, surpassing the state-of-the-art, with over 200 times higher throughput compared to traditional SMO methods.

Index Terms—Deep learning, design automation, design for manufacture.

I. INTRODUCTION

OPTICAL lithography plays a vital role in integrated circuit (IC) manufacturing, using light to selectively expose a photosensitive material and meticulously pattern IC’s thin film layers. State-of-the-art (SOTA) resolution enhancement techniques like source mask optimization (SMO) and multiple patterning have facilitated the creation of complex ICs, characterized by smaller feature sizes and enhanced transistor densities. In essence, source optimization (SO) adjusts the incidence angles or intensity distributions of light rays emitted from the illumination system. Conversely, mask optimization (MO) regulates the amplitude of the transmitted light rays via mask, as shown in Fig. 1. Compared to MO, which is more commonly known as optical proximity correction (OPC), SMO can gain a larger manufacturing process window by adjusting the source intensity distribution. This feature has generated significant interest among semiconductor foundries, inciting extensive research to prolong the viability of 193-nm optical lithography [1], [2]. The advent of freeform illumination has enabled pixelated SMO techniques to enhance optimization freedom degrees, consequently boosting the imaging performance of lithography systems [3], [4]. However, the intricate source representation significantly affects the pixelated SMO method’s performance. Additionally, the escalating density of IC layouts triggers a surge in the data volume processed by SMO algorithms. Pixel-based optimization introduces numerous optimization variables, thereby escalating the computational complexity of the optimization process. Collectively, these factors impose substantial challenges on the computational efficiency of SMO algorithms.

The advent of Hopkins’ approach [5] and its simplified formulation, the “sum of coherent systems” (SOCS) [6], have triggered a significant reduction in MO computational complexity via singular value decomposition (SVD). Consequently, research has veered toward the exploration of distinct acceleration techniques for SO and MO [7], [8], [9], [10].

Recent years have seen a notable surge in research dedicated to MO acceleration [11], [12], [13], [14] using GPU and deep

The illumination system of advanced lithography machines has made significant progress, achieving what is called freeform illumination. This intensity distribution of source map is composed of many pixelated illuminators, and the intensity of each pixel can be controlled by computer programs. The pixelated illumination system is necessary for SO. However, SO, which is dictated by Abbe’s imaging approach, presents a challenge for the direct application of GPU acceleration techniques. Wang et al. [16] proposed to use compressive sensing to reduce the computational complexity of the SMO algorithm. ICC-CPU [17] puts forth a sampling-based imaging model employing heuristic algorithms to augment SMO efficiency and performance. SoulNet [18] utilizes an autoencoder to learn the model-based source correlating with each layout. Despite this, SoulNet falls short in accomplishing iterative SMO and offers limited optimization capabilities for the overall result. At present, while MO optimization time has been reduced to seconds for the ICCAD13 benchmark [15], CPU-based SO still requires several hours of optimization, creating a substantial bottleneck within the overall SMO workflow.

To navigate these challenges and foster advancements in SMO research, we introduce the Diff-SMO framework, composed of three key components. Initially, we have constructed a lithography algorithm tailored for GPU parallel acceleration, exploiting the intrinsic properties of the rigorous Abbe model. Subsequently, building on previous accomplishments in GPU/DNN acceleration for MO, we have achieved GPU acceleration of the SO algorithm (GPU-SO), culminating in full GPU acceleration of the entire SMO workflow. Lastly, we present Diff-SO, which incorporates a conditional discrete diffusion model for quick production of mask-aware, near-optimal sources. The diffusion process imitates the progressive turning on/off of the source in a discrete state space. The resultant solution of diffusion model provides an improved initial choice for GPU-SO to facilitate further precise optimization. By reducing the number of iterations required by GPU-SO, Diff-SO expedites the SMO process and enhances the optimization results. Our contributions can be summarized as follows.

1) We present the first academic effort, to our knowledge, to design a GPU-accelerated lithography simulator based on Abbe’s aerial imaging theory.
2) We devise GPU-SO, a GPU acceleration algorithm specifically for SO, thus facilitating the comprehensive integration of GPUs across the entire SMO process.
3) We propose Diff-SO, an innovative conditional discrete diffusion model that generates near-optimal sources, enhancing the efficiency and performance of SMO significantly.
4) Our experimental data showcase the effectiveness of our method, demonstrating an average improvement of 9% over SOTA MO techniques. Furthermore, in comparison to prevalent CPU-based SOTA SMO methods, Diff-SMO provides an 8% performance increase along with a remarkable speed boost of nearly 225 times.

II. PRELIMINARIES

In this section, we will discuss the fundamental concepts of lithography imaging theory. We will specifically focus on Abbe’s method, which serves as the foundation for SO algorithms. Additionally, we will explore Hopkins’ method, a widely used approach for MO algorithms. Furthermore, we will define the SMO problem and delve into the evaluation metrics associated with it.

A. Lithography Simulation

In a typical lithography simulation, the lithographic imaging model comprises two essential components: 1) the illumination source and 2) the projector system, as depicted in Figs. 1 and 2. During the manufacturing process, the light rays propagate through the mask and produce diffracted light with feature pattern information. The intensity distribution of aerial image can be formulated via lithography theory [4] as

$$I(x, y) = \iint \int \int_{-\infty}^{\infty} J(f, g) \cdot F(M)(f', g') \cdot F(M)^* (f'', g'') \cdot H(f + f', g + g') \cdot \exp(-j2\pi ((f' - f')x + (g' - g')y)) \, df \, dg \, df' \, dg' \, df'' \, dg''$$

where $I$ is the aerial intensity distribution; $J$ is the lithographic illumination source; $H$ denotes the optical transfer function (OTF) of the projector system; $F(M)$ represents the frequency spectrum of the mask pattern $M$ obtained via $2 - D$ fast Fourier transform ($FFT$) $F(\cdot)$; $^*$ is the Hermitian transpose operator; $(x, y)$ denotes the spatial coordinates on the image plane; $(f, g)$, $(f', g')$ and $(f'', g'')$ represent the normalized-frequency-domain coordinates of the pupil, mask spectrum and its conjugate, respectively. Two formulas are used to calculate the partially coherent aerial image formation in (1).
As illustrated in Fig. 3, one is Abbe’s approach [19] and the other is Hopkins’ method [5].

Abbe’s Approach: In Abbe’s approach, the source illuminator is discretized, as illustrated in Fig. 3(a). The contribution from each source point toward the final aerial image is calculated independently and then summed. For convenience, we introduce the illumination cross-coefficients (ICC) as

\[
ICC(x, y; f, g) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} H(f + f', g + g') F(M)(f', g') \exp(-j2\pi(f'x + g'y)) df' dg'.
\]  

(2)

Equation (1) can be formulated in Abbe’s approach

\[
I(x, y) = \int_{-\infty}^{\infty} J(f, g) ICC(x, y; f, g) df' dg'.
\]  

(3)

Hopkins’ Approach: As depicted in Fig. 3(b), Hopkins’ approach separating the illuminator and projector system from mask has the following expression that is:

\[
I(x, y) = \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} T(f', g'; f'', g'') F(M)(f', g') F(M)**(f'', g'') \exp(-j2\pi((f' - f'')x + (g' - g'')y)) df' dg' df'' dg''.
\]  

(4)

where \( T \) is the transmission cross-coefficients TCC given by

\[
TCC(f', g'; f'', g'') = \int_{-\infty}^{\infty} J(f, g) H(f + f', g + g') H^*(f + f'', g + g'') df' dg'.
\]  

(5)

To reduce the calculation complexity, an approximation solution for the Hopkins imaging equations called Sum of coherent source (SOCS) decomposes the TCC spectrum by applying SVD as follows:

\[
TCC(f', g'; f'', g'') = \sum_{q=1}^{\infty} \kappa_q \Phi_q(f', g') \Phi_q^*(f'', g'')
\]  

(6)

where \( \kappa_q \) and \( \Phi_q \) are \( q \)’th eigenvalue and eigenvector of TCC. Because the eigenvalue \( \kappa_q \) is rapidly decay with increasing \( q \), we can keep the \( Q \) largest eigenvalues for faster calculation. Substituting (6) into (4) then performing Inverse Fast Fourier Transform (IFFT), the SOCS can be expressed to spatial position via (7) that is

\[
I(x, y) = \sum_{q=1}^{Q} \kappa_q \Phi_q(x, y) \otimes M(x, y)
\]  

(7)

where \( \Phi_q(x, y) \) and \( M(x, y) \) are the spatial distribution of \( \Phi_q \) and \( F(M) \), respectively. \( \otimes \) and \(| \cdot |\) are the convolution and absolute operator.

Hopkins’ Approach Versus Abbe’s Approach: Hopkins’ method (7) applies the SOCS to decrease computational complexity from \( O(n^6) \) to \( O(Q \times n^4) \), considering \( M \in \mathbb{R}^{n \times n} \). Importantly, it separates mask from optical system, rendering Hopkins’ method more suitable for simulating images while optimizing masks. With constant optical imaging conditions, Hopkins’ method outperforms Abbe’s in runtime, making it a favored choice in MO algorithms [7, 11, 12, 13, 14]. In contrast, Abbe’s method in (3), sums up the impact of all source points to generate the final aerial image. This method is inherently suited for SO due to the source’s discretization. Consequently, to achieve SO and SMO, devising a lithography simulator founded on Abbe’s method is a prerequisite. The comparison of Hopkins’ and Abbe’s imaging is illustrated in Fig. 3.

**B. Evaluation Metrics**

**Definition 1 (Squared L2 Error):** Given target pattern \( Z_t \) and wafer pattern under nominal process condition \( Z \), the squared L2 error is calculated as \( ||Z - Z_t||_2^2 \).

**Definition 2 (Process Variation Band (PVB)):** PVB [15] is used in manufacturing to represent the expected range of variation in a production process. PVB denotes the XOR area between the aerial images \( Z_{in} \) and \( Z_{out} \) under the min and max lithography process conditions.

**Definition 3 (Edge Placement Error (EPE)):** EPE [15] refers to the deviation between the intended position of a feature on a wafer and its actual position after lithography.

**C. Problem Formulation**

To formulate the problem of SMO, we first introduce the necessary terminologies. Given a target image represented as \( Z_t \), we define a function \( F \) which represents the weighted sum of three components: 1) the L2 error; 2) PVB; and 3) EPE that arise after the lithography process. Under given circumstances: 1) if the source parameters \( J \) and the projector system are predetermined, the task of MO is to discern the ideal mask configuration that minimizes \( F \) and 2) conversely, with a fixed mask in place, SO focuses on adjusting the source \( J \) to achieve a minimal \( F \). The overarching goal of SMO is to concurrently identify the optimal configurations for both the source, denoted as \( \hat{J} \) and the mask, denoted as \( \hat{M} \), to minimize \( F \). This can be mathematically represented as

\[
(\hat{J}, \hat{M}) = \arg \min_{(J,M)} F(J, M).
\]  

(8)
III. ALGORITHM

This section outlines the Diff-SMO framework, as visualized in Fig. 4. Initially, we address the GPU-accelerated SO algorithm in Section III-A. Then, Section III-B introduces the diffusion model, crucial for mask-aware source creation. Finally, Section III-C delves into the entire Diff-SMO process. We simplify our terminology as follows: GPU-SO for the GPU-accelerated SO algorithm, Diff-SO for the application of diffusion in source generation and its subsequent optimization with GPU-SO, and Diff-SMO for the combined framework of Diff-SO and the MO algorithm, as depicted in Fig. 4.

A. GPU-Accelerated SO Algorithm

Computation Complexity Analysis of Vanilla SO: In numerical terms, (3) must be resolved for a discrete number of mutually independent source points as follows:

$$I(x_i, y_i) = \sum_{f=1}^{N_s} \sum_{g=1}^{N_s} J(f, g) \text{ICC}(x_i, y_i; f, g) \tag{9}$$

where $N_s$ denotes the lateral dimension’s pixelated source sampling number. We designate the $S$ operator as the reshape operation

$$\mathbb{R}^{N_s \times N_s} \xrightarrow{S} \mathbb{R}^{N_s \times N_s} \xrightarrow{S^{-1}} \mathbb{R}^{N_s \times N_s}. \tag{10}$$

As a result, (9) can be translated into (11) via matrix multiplication

$$I = S^{-1}(\text{ICC} S(J)) \tag{11}$$

where $I \in \mathbb{R}^{N_s \times N_s}$, $\text{ICC} \in \mathbb{R}^{N_s^2 \times N_s^2}$, and $J \in \mathbb{N}^{N_s \times N_s}$, respectively, and $N$ is the mask spatial domain’s sampling number. The computational complexity is $O(n^6)$. In our implementation, specifically, $N = 2048$ and $N_s$ ranges from 35 to 58. The computation of large matrix products, between ICC and $J$, incurs significant time cost and necessitates careful memory allocation. This computational demand serves as a substantial impediment to the execution time of the forward aerial imaging process. Furthermore, the situation worsens during gradient calculation, where the process often necessitates a duration several magnitudes greater than the forward computation time.

GPU-Accelerated Abbe Aerial Imaging: Historically, the literature has largely depended on acceleration techniques employing compressive sensing or heuristic approaches [17] to selectively observe certain points from the extensive ICC matrix. These strategies often tradeoff a degree of accuracy for computation speed enhancement. However, our approach distinguishes itself by leveraging GPU acceleration to fundamentally address the slow computation issue inherent in the Abbe model. Importantly, we recognize that each discrete light source computation in the Abbe model, as per (9), is independent and noninterfering. This characteristic renders the Abbe model particularly amenable to parallel computing acceleration via GPUs.

The algorithm outlined in Algorithm 1 presents a new approach to GPU-accelerated aerial imaging generation based on Abbe’s approach. It starts with noise elimination and the selective retention of source points that contribute to the generation process, as indicated in line 2. By applying a threshold, we can reduce computational overhead caused by zero-valued sources. The method contrasts with prior techniques which relied on fixed-dimensional matrix multiplication. Instead, our algorithm allows for finer acceleration granularity at the
The contributions of all source points are summed up to generate the aerial image. Each discretized source is assigned to an individual thread for computation. This is achieved by allowing independent filtering and computation for each source point, as detailed in line 5 and Fig. 5. Building on Abbe’s imaging theory, the projector transfer function is computed as per line 7, and can be visualized as a low-pass filter of the spatial frequency domain, given by (12)

\[
H(f, g) = \begin{cases} 
1, & \sqrt{f^2 + g^2} \leq \text{NA}/\lambda \\
0, & \text{otherwise}
\end{cases}
\]

(12)

Similarly, \(J(f, g)\) is expressed by (13), where \(\text{NA}/\lambda\) is the cut-off spatial frequency defined by the wavelength \(\lambda\) and numerical aperture \(\text{NA}\) of the projection system.

\[
J(f, g) = \begin{cases} 
J(f, g), & \sqrt{f^2 + g^2} \leq \sigma \text{NA}/\lambda \\
0, & \text{otherwise}
\end{cases}
\]

(13)

Further, the partial coherence factor \(\sigma\) is determined by the size ratio of the source image and the pupil. \(\sigma \text{NA}/\lambda\) represents the maximum spatial working region of the source, relative to the spatial frequency space of the image domain, as evaluated in line 8. By applying the IFFT to the result of an element-wise product (\(\odot\)) of \(H_s\) and \(F(M_s)\), as shown in line 9, we can compute the contribution from each point. The final aerial image \(I\) is obtained by summing up these individual contributions from each source point, each calculated within a separate thread, as summarized in line 12.

**GPU-Accelerated SO**: The GPU-accelerated SO algorithm, shown in Algorithm 2, employs freeform illumination in the pixellated SMO process, which grids the source pattern \(J\) into \(N_s \times N_s\) pixels, as is depicted in Fig. 5. Each pixel’s value signifies the intensity of the associated source point, falling within the range of \([0, 1]\). To ensure the SMO flow’s differentiability, we introduce the source parameters \(J_\theta\), as initialized by

\[
J_\theta(f, g) = \begin{cases} 
1, & J_0(f, g) = 1 \\
-1, & J_0(f, g) = 0
\end{cases}
\]

(14)

where \(J_\theta \in \mathbb{R}^{N_s \times N_s}\) and can hold any real value in the range \((-\infty, \infty)\). Meanwhile, \(J_0\) represents the initial pixellated source pattern derived from the parametric representations of source templates. Some examples of \(J_0\) are visualized in Fig. 6. During the optimization process, source \(J\) is obtained by applying the Sigmoid function to the parameters \(J_\theta\)

\[
J = \text{Sigmoid}(\alpha \cdot J_\theta) = \frac{1}{1 + \exp(-\alpha \cdot J_\theta)}
\]

(15)

where \(\alpha\) is a hyperparameter to define the steepness of the Sigmoid function. Afterward, the aerial image \(I\) is calculated using the Abbe_litho function (introduced in Algorithm 1), with \(J\) and mask \(M\) as inputs. The Algorithm 2 utilizes a simple resist model in line 5, which applies the Sigmoid function to \(I\)

\[
Z = \frac{1}{1 + \exp(-\beta \cdot (I - I_{tr}))}
\]

(16)

where \(Z\) is the resist pattern; \(I_{tr}\) is the intensity threshold; \(\beta\) is a hyperparameter for steepness. In line 6, the mean-squared error \(L\) is then computed between the predicted resist image \(Z\) and the target \(Z_t\), where \(L = \|Z - Z_t\|^2\). The source parameters \(J_\theta\) are subsequently updated in line 7 via gradient backpropagation using the Adam optimizer. The optimization process persists until the change in loss \(\Delta L\) between successive iterations falls below a threshold \(\delta\).

**GPU-Accelerated SO Versus Multicore CPU-Accelerated SO**: Central processing units (CPUs) act as the primary computational elements, executing essential logical, arithmetic, and input/output operations. They also manage the distribution of tasks to different components and subsystems in a computer system. In modern times, CPUs are frequently designed with a multicore structure, which means multiple processing units are housed within a single IC. This design approach not only conserves power but also improves performance, enabling efficient parallel processing of simultaneous tasks.

The previously mentioned GPU-accelerated SO’s parallel algorithm can be fully migrated to a multicore CPU. We have also implemented a CPU parallel acceleration algorithm based on multiprocessing. However, compared to the multicore CPU-accelerated SO, the GPU-accelerated SO still has the following advantages. First, because the entire SO optimization process doesn’t require extensive IO reading operations, it can utilize GPU pinned memory to speed up computations. Additionally, GPUs offer greater memory bandwidth, saving significant time during memory access. Furthermore, for large kernel size operations, GPUs can compute FFT/IFFT more quickly, which is frequently invoked throughout the SO process. Also, previous SOA MO researches have shown that GPU-accelerated MO is effective. A comprehensive GPU-accelerated SMO algorithm can greatly reduce the time wasted transferring data between the CPU and GPU. Lastly, multicore CPUs might be more expensive than GPUs. For tasks like SO that require massive parallelism, GPUs are not only more suitable but also more cost-effective than CPUs.

**B. Diffusion Models for Mask-Aware Source Generation**

Currently, the estimation of our GPU-accelerated SO is initiated from a template source as defined in (14) and illustrated in Fig. 6. Intuitively, a mask-aware initialization can reduce the number of iterations and accelerate the optimization process. The task of exact source initialization for our GPU-accelerated SO is difficult to formulate using traditional rule-based methodology. However, an approximate initialization is sufficient in this case. Therefore, we resort to learning-based methods to solve this problem.
The source initialization can be formulated as a conditional binary discriminative task: given the mask and an empty source matrix, the model learns to determine which pixels in the source matrix should be turned on/off. Using a discriminative model to judge the state of every source pixel can be a natural choice but the large amount of discriminative heads can be too heavy. Instead of that, we introduce an efficient generative model, namely, the conditional discrete diffusion model, in our method to generate source matrix from given condition.

**Conditional Discrete Diffusion Denoising Model**: Typical generative models, such as GANs and diffusion model [20], are commonly applied in the continuous value space, where each generated pixel is represented as a floating-point number within a predefined range. However, in our source initialization case, the state of each pixel in the source matrix is binary (on/off). Applying a threshold to a continuous pixel value and converting it into a binary value may result in information truncation and compromise performance. In contrast to the traditional diffusion models used in computer vision, we have made several important customizations to enhance the diffusion model and directly synthesize a discrete source matrix from the given mask. To derive these enhancements, we first reformulate the problem.

In our approach, we employ a forward process to simulate the gradual turning off of the light in the given source matrix. Similarly, we utilize a reverse process to mimic the gradual turning on process. Our conditional discrete diffusion model defines a Markov chain to represent both the forward turn-off process and the reverse turn-on process, as illustrated in Fig. 7. The lengths of both processes are denoted as $K$. During the $k$th step of the forward process, $x_k \in \{0, 1\}$ represents an entry in the source matrix $J$. The forward step can be described by the equation

$$q(x_k | x_{k-1}) := \text{Cat}(x_k; p = x_{k-1} Q_k)$$

where $x_k$ is the one-hot version of the entry $x_k$, $\text{Cat}(x; p)$ is a categorical distribution over the row vector $x$ with probabilities given by the row vector $p$, $x_{k-1} Q_k$ can be understood as a row vector-matrix product and the transition probability matrix $[Q_k]_{ij} = q(x_k = j | x_{k-1} = i)$ is defined to describe the state transition probability for each $x$ at the $k$th step. The $Q_k$ is applied to each entry in the source matrix independently and $q$ factorizes over these higher dimensions as well. On the other side, in the reverse process, given the mask $M$, the neural network aims to predict the conditional categorical distribution probability $p_0(x_{k-1} | x_k, M)$ over each pixel to recover the original source matrix.

As we described before, in the forward turn-off process, the distribution of every entry $x_k$ becomes large. So given any $x_0$ and mask $M$, the distribution of every entry $x_k$ should follows:

$$q(x_k | x_0, M) \rightarrow [1, 0], \text{ when } k \rightarrow K.$$  

Therefore, we design a transition matrix $Q_k$ with an absorbing state 0 for the turn off process

$$Q_k = \begin{bmatrix} 1 & 0 \\ \beta_k & 1 - \beta_k \end{bmatrix}$$

where $\beta_k \in (0, 1)$ is the hyperparameter controlling the flipping probability. In order to ensure that at the $k$th step, $(k/K)$ information about given source matrix is lost, following the suggestion of [21], we use an increasing noise schedule for $\beta_k$:

$$\beta_k = \frac{1}{K - k + 1}, \quad k = 1, \ldots, K.$$  

**Training Diffusion Model**: To training the conditional discrete diffusion model for source matrix generation, the training objective at step $k$ is to minimize the loss function

$$L = D_{KL}(q(x_k | x_{k-1}, x_0) || p_0(x_{k-1} | x_k, M)) - \eta \log p_0(x_k | x_0, M)$$

where $\eta$ is a hyperparameter to balance the loss terms.

Given a source matrix $J_0$, we randomly sample a target step $k$ from an uniform distribution that is defined from 1 to $K$ first, and expect to get a noisy sample $J_k$. Fortunately, we can explicitly derive that $x_k$ obeys the following categorical distribution:

$$q(x_k | x_0) = \text{Cat}(x_k; p = x_0 Q_k)$$

The choice of transition probability matrix $Q_k$ is critical. As we described before, in the forward turn-off process, the stationary state of every entry in the source matrix should be off (0 in our case) when $k$ becomes large. So given any $x_0$ and mask $M$, the distribution of every entry $x_k$ should follows:

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**Algorithm 3 Training of Discrete Diffusion Model**

**Input**: Ground Truth $\hat{J}$; Source template $J_0$; Mask $M$.

**Output**: Well-trained model $p_0$.

1: repeat
2: $k \leftarrow$ Random$(1, K)$;
3: $J_k \leftarrow \text{Turn_off}(J, k); \quad \triangleright \text{Turn_off, Eq. (22)}$
4: $J_0 \leftarrow p_0(J_k, J_0, M, k)$;
5: $L \leftarrow \text{Loss}(J_0, J_k); \quad \triangleright \text{Eq. (21)}$
6: $p_0 \leftarrow \text{Update}(p_0, L)$;
7: until Finished
8: return Well-trained model $p_0$;

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4: $J_0 \leftarrow p_0(J_k, J_0, M, k)$;
5: $L \leftarrow \text{Loss}(J_0, J_k); \quad \triangleright \text{Eq. (21)}$
6: $p_0 \leftarrow \text{Update}(p_0, L)$;
7: until Finished
8: return Well-trained model $p_0$;
Algorithm 4 Generating Mask-Aware Source

**Input:** Source template \( J_0 \); Mask \( M \); Model \( p_0 \).

**Output:** Estimated mask-aware source \( \hat{J}_0 \).

1. \( J_k \leftarrow \text{all}_0 \text{matrix} \).
2. **for** \( k \in \{1, \ldots, K\} \) **do**
3. \( \hat{J}_0 \leftarrow p_0(J_0, J_0, M, k); \)
4. \( J_k \leftarrow J_{k-1}(\hat{J}_0, J_k); \quad \triangleright \text{Turn_on}, \text{Eq.} \ (23) \)
5. **end for**
6. **return** Estimated mask-aware source \( \hat{J}_0 \);

where \( \hat{O}_k = Q_1 Q_2 \ldots Q_k \). Then, we can directly sample from the above distribution to obtain \( J_k \), instead of turning off \( k \) times.

After sampling \( J_k \), we extract the real \( \hat{M} \) and imaginary \( \tilde{M} \) parts from the spectrum of given mask \( M \). Then, we concatenate \( J_k, \hat{M}, \tilde{M} \) and the template source together and feed them into the neural network with the embedding of the time step \( k \). The template source is defined in (14). The neural network will predict the logit of the posterior \( \text{logits} \) of the posterior time step \( k \). The template source is calculated as following:

\[
p_\theta(x_{k-1}|x_k, M) = \sum_{x_0} q(x_{k-1}|x_k, x_0) p_\theta(x_0|x_k, M) \tag{23}
\]

where the term \( x_0 \) will visit every possible state of \( x_0 \). And \( q(x_{k-1}|x_k, x_0) \) has a closed form according to (17) and Bayes’ theorem

\[
q(x_{k-1}|x_k, x_0) = \text{Cat} \left( x_{k-1}; p = \frac{x_k Q_k^\top \odot x_0 Q_{k-1}^\top}{x_k Q_k x_k^\top} \right) \tag{24}
\]

where \( \odot \) is a pixel-wise multiplication.

So far, all items in the loss function have been obtained, and the neural network can be trained by the commonly used gradient descent method. The training process is summarized in Algorithm 3.

Generating Mask-Aware Source Matrix: Once the model is well-trained, given the mask \( M \) and the type of template source, we can extract a mask-aware source matrix by starting from a all-zero (dark) source matrix and gradually turn on the pixel with the reverse procedure. The generating process can be expressed by

\[
p_\theta(\hat{J}_0|J_K, M) = p_\theta(\hat{J}_0|J_1, M) \prod_{k=2}^{K} p_\theta(J_{k-1}|J_k, M) \tag{25}
\]

where \( J_0 \) is the estimated source matrix at step \( k \) and \( \hat{J}_0 \) is the newly sampled source matrix. When generation finished, the generated source matrix \( \hat{J}_0 \) is naturally a binary one where each entry equals either zero or one. And the generated source matrix will be utilized as the initialization of our GPU-accelerated SO. The inference process is summarized in Algorithm 4.

C. Diff-SMO Flow

As illustrated in Fig. 4, the process commences with the loading of the source \( J_0 \) and mask \( M \). By using the source template and mask spectrum as conditional factors, a well-trained diffusion model is deployed to efficiently generate a near-optimal source approximation, \( \hat{J}_0 \). This approximation is then fed into our GPU-accelerated SO solver for a more precise optimization, leading to the optimal source \( J \). When the predetermined stopping criteria for SO are satisfied, a TCC matrix will be calculated by integrating the optimized source with the pupil data following (5). Subsequently, the TCC matrix is utilized in MO solvers, as referenced in [7], [12], [13] and [14], to further optimize the mask. The refined mask and source are iteratively fed back into the Diff-SMO optimizer, guiding the optimization process until convergence is achieved for both the source and mask components.

IV. EXPERIMENTS

The Diff-SMO framework, developed in the Pytorch framework, was tested on a Linux System using a single Nvidia RTX 3090 GPU card. The multicore CPU-accelerated SO algorithm runs on AMD EPYC 9554P with 64 cores using the 7-nm Zen architecture. This architecture supports up to 64 cores and 128 threads. Hyperthreading and Multithreading technologies allow each core to have two virtual cores, enhancing performance. We conducted a calibration of the lithographic system to match the simulator [15] used in ICCAD 2013 before starting our experiments. The hyperparameter settings are as follows: \( \sigma = 1; \alpha = 8; \beta = 30; I_r = 0.225; \delta = 10^{-4}; \lambda = 193 \text{~nm}; \eta = 0.001; K = 100; \text{inner radius} \sigma_1 = 0.63; \text{outer radius} \sigma_3 = 0.95; \) and refractive index is 1.414. For simplicity, we denote squared L2 error, PVB, and EPE as “L2,” “PVB,” and “EPE,” respectively.

We evaluate the performance of the Diff-SMO framework on the ICCAD 2013 CAD contest benchmark [15] using ten 4 \( \mu \text{m}^2 \) tiles of 32-nm industrial layouts. The training set for the diffusion model is obtained from [11], which contains 4K \( 4 \mu \text{m}^2 \) generated tiles following the same design rules as in [15]. The optimal ground-truth source for each mask is obtained by fully optimizing the template source with GPU-SO, using the corresponding mask as input.

Following [20], [22], the U-Net is selected as our backbone for diffusion model. The model predicts the posterior distribution during the reverse diffusion process. All input images are resized to \( 32 \times 32 \) and consists of four feature map resolutions: [32 \( \times \) 32, 16 \( \times \) 16, 8 \( \times \) 8, 4 \( \times \) 4]. Each resolution level contains two convolutional residual blocks, with the number of convolution channels being [128, 256, 256, 256] correspondingly. A self-attention block is placed at the 16 \( \times \) 16 resolution level. Additionally, the time step \( k \) is incorporated into each residual block through the sinusoidal position embedding [23].

A. Result Comparison With the SOTA

In this section, we will first validate the efficiency of Diff-SO. Following that, we will compare Diff-SO with the previous SOTA SO methods. The approach is to use different MO methods as a foundation, combined with various SO methods, to validate the performance of different SMO methods. Then, we compared the EPE and Runtime performance of different SMO combinations. Finally, we compared the
TABLE I
EFFECTIVENESS OF DIFF-SO

<table>
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</table>

L2 and PVB unit: nm². The ratio calculates the proportion of “w/ Diff-SO” and “w/o SO” of the corresponding MO algorithm.

Diff-SO Effectiveness Evaluation: We conduct performance evaluations of Diff-SO with various MO algorithms [7], [12], [13], [14] in Table I and observe that combining any MO algorithm with the Diff-SO algorithm consistently improves L2 and PVB results. It is worth noting that in Table I, the “ratio” row represents the percentages calculated between results obtained using only MO (labeled as “w/o SO”) and the results achieved by incorporating Diff-SO with the corresponding MO method within the Diff-SMO framework (labeled as “w/ Diff-SO”). When the vanilla MO algorithm MOSAIC [7] is combined with the Diff-SO algorithm, a significant reduction of 18% in L2 error and 7% in PVB is achieved. The SOTA level-set-based MO algorithm, DevelSet [12], can reduce the L2 error by 12.5% and the PVB by 7.1% when combined with Diff-SO. When combined with A2ILT [13], the SOTA pixel-based MO algorithm, Diff-SO achieves a reduction of 11.8% in L2 error and 8.1% in PVB, respectively. The latest MO research, Multi-ILT [14], has made groundbreaking progress by utilizing multilevel ILT combined with SRAFs. However, Diff-SO still manages to achieve an overall reduction of 9.5% in L2 error and 7.8% in PVB. These experiments validate the effectiveness of Diff-SO within the Diff-SMO framework, demonstrating its capability to yield an average enhancement of approximately 13% in L2 error and 7.5% in PVB across various SOTA MO methods.

Diff-SMO Result Comparison With SOTA SMO Methods: As depicted in Figs. 8 and 9, the x-axis represents different MO algorithms. Each point represents a combination of different SMO algorithms. We compare the performance of Diff-SO with previous SOTA SO solvers, including ICC-CPU [17] and SoulNet [18]. Regarding MO solvers, unlike the original MOSAIC, our implementation entails GPU acceleration to enhance MOSAIC’s computational efficiency. From Figs. 8 and 9, it can be observed that our approach, including GPU-SO and Diff-SO, outperforms the previous SOTA SO methods in terms of both L2 and PVB metrics. Specifically, for Diff-SO, the L2 metric shows an average improvement of 9.2% and 14.3% compared to ICC-CPU and SoulNet, respectively. Regarding the PVB metric, Diff-SO demonstrates an average enhancement of 4.7% and 6.9% over ICC-CPU and SoulNet, respectively. It is noteworthy that the magnitude of MO optimization objectives directly influences the potential optimization space for Diff-SO. For example, when Multi-ILT is chosen as the MO solver, Diff-SO exhibits a modest 9.5% improvement in L2 alone. However, employing a less...
sophisticated MO solver, such as MOSAIC, allows Diff-SO to achieve a significant 18% enhancement in L2. Several sample results of Diff-SMO are depicted in Fig. 10, utilizing an annular source as the initial source condition.

**EPE and Runtime Comparison of SMO Algorithms:** In Tables II and III, we further compare the performance of EPE and runtime. The “average” in Table II for the row MO-solver represents the average results obtained by testing MOSAIC, DevelSet, A2ILT, and Multi-ILT as MO solvers. The table presents the average outcomes of these MO solvers when different SO solvers are used. Diff-SO achieves a reduction of 1.2% and 3.1% in EPE compared to ICC-CPU and SoulNet, respectively, as shown in Table II. Table III presents a comprehensive runtime comparison of the SMO algorithm execution, encompassing both the SO and MO phases. SoulNet [18] directly generates the source using an autoencoder, omitting the iterative SMO optimization process. Consequently, this leads to a modest 0.3-s increase in runtime for Diff-SMO compared to SoulNet. The observed increase can be attributed to the subsequent feeding of the generated source from the diffusion model into GPU-SO for further fine-grained optimization steps. However, the design has proven to be valuable, as the 0.3-s increase in runtime has resulted in a notable 14.3% decrease in L2 error. Moreover, in comparison to the conventional SO algorithm, ICC-CPU [17], Diff-SMO demonstrates a remarkable speed enhancement, achieving a nearly 225-fold improvement in the overall SMO runtime while producing better results.

**SoulNet Versus Diffusion Model:** In Table IV, we compare our diffusion model with SOTA learning-based method SoulNet [18]. The term “mIoU” stands for Mean Intersection over Union, which represents the average ratio of intersection over union. “pixelAcc” represents pixel accuracy, defined as the percentage of pixels correctly classified in an image. Our diffusion model exhibits improvements of 4.79% in mIoU and 2.44% in pixelAcc. Additionally, it requires less inference time compared to SoulNet due to its simpler network architecture.

**B. Ablation Study**

**CPU-SO Versus GPU-SO:** We conducted a series of experiments to validate the efficiency of the GPU-accelerated SO algorithm. First, we compared it with the acceleration of SO using a multicore CPU with multithreads, referred to as CPU-SO. The CPU we used has 64 cores, supporting up to 128 threads. By varying the number of threads used (from 1 to 128), we averaged the CPU runtime for ten test cases, and the results are plotted in Fig. 11. The detailed experimental results for each sampled thread are shown in Table V. From Fig. 11 and Table V, it can be observed that between 1 to 32 threads, the runtime shows a significant decrease as the number of threads increases. When the number of threads exceeds 32, the runtime remains relatively stable. At 64 threads,
which is precisely the number of CPU cores, the best-runtime performance is achieved, taking 65 s to complete SO. In contrast, the GPU’s runtime is 16 s. Notably, when the number of threads exceeds 64, even at 65, there’s a noticeable increase in runtime. We believe this is due to CPU oversubscription, causing different processes to compete for CPU resources, resulting in reduced efficiency.

We also conducted another set of experiments. Given that the best-experimental results are achieved with 64 threads, how would the runtime for both the 64 threads CPU and GPU be affected if we change the number of source points that need to be optimized in each iteration, thereby altering the overall computational load. The experimental results are shown in Fig. 12. Observing (9), when ICC is fixed, the total computational load of SO is directly proportional to the source points of the illumination source J. By adjusting the size of the source area, we controlled the number of source points to be computed, ranging from 20 to 2500. Now we can compare the runtime performance under varying computational loads to observe how CPU-SO and GPU-SO perform as the amount of computation changes. As seen in Fig. 12(a), as the number of effective source points increases, resulting in a larger computational load, the runtime required for SO increases. However, it’s noteworthy that regardless of whether the number of source points is greater or less than 64, the ratio of runtime between the 64-thread CPU and GPU consistently remains around 4 times, as depicted in Fig. 12(b).

In summary, as shown in Fig. 11, when the computational workload is constant, we observe that the total runtime dramatically decreases with the increase in parallel threads, reaching a minimum at 64 threads. The runtime is approximately four times that of GPU-SO. Then, in Fig. 12, we vary the computational workload and compare the 64-threads CPU-based SO with the GPU-SO, maintaining a runtime that is roughly four times longer. This confirms the necessity and efficiency of our GPU-accelerated SO algorithm.

**GPU-SO Versus Diff-SO:** In Table III, we present a comparative analysis of runtime between GPU-SO and Diff-SO under the integration with various MO methods. It is observable that in the absence of the MO process, GPU-SO, on average, takes an additional 14.7 s to execute, approximately 12.6 times that of Diff-SO. By utilizing a diffusion model to generate near-optimal sources as initial inputs for GPU-SO iterations, Diff-SO significantly curtails the SO time by 92%, leading to enhanced optimization outcomes. From Figs. 8 and 9, it is evident that compared to GPU-SO, Diff-SO also augments the overall SMO performance, reducing the L2 error by 2% and the PVB by 1.5%. Relative to MO, the diffusion model is exceptionally suitable for SO. Typically, a mask clip handled is of size $2048 \times 2048$ pixels, and employing the diffusion model for such dimensions necessitates substantial memory and entails extensive computation time. In contrast, common source sizes range from $35 \times 35$ to $60 \times 60$. For sources of these dimensions, a modest GPU can leverage the diffusion model to generate sources with remarkable precision.

**V. CONCLUSION**

We introduce Diff-SMO, an advanced and accelerated solution for SMO. Our framework innovatively expands on existing MO algorithms, with a primary focus on accelerating and refining SO. Our distinctive contribution is the formulation of a GPU-accelerated algorithm grounded in Abbe’s imaging theory, which enables full GPU acceleration throughout the SMO flow. Moreover, we integrate a novel diffusion model that gradually learns to activate or deactivate source points, resulting in rapid generation of approximately optimal solutions, significantly boosting the SMO process. Experimental outcomes validate the superiority of Diff-SMO over combined prior SOTA SO and MO algorithms. Remarkably, it accomplishes an acceleration exceeding 200 times compared to conventional SMO methods, concurrently producing superior results. In future work, we plan to open-source our lithography model and SMO framework, fostering further research in this area.

**REFERENCES**


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