# A stochastic gradient descent method for computational design of random rough surfaces in solar cells 

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#### Abstract

In this work, we develop a stochastic gradient descent method for the computational optimal design of random rough surfaces in thin-film solar cells. We formulate the design problems as random PDE-constrained optimization problems and seek the optimal statistical parameters for the random surfaces. The optimizations at fixed frequency as well as at multiple frequencies and multiple incident angles are investigated. To evaluate the gradient of the objective function, we derive the shape derivatives for the interfaces and apply the adjoint state method to perform the computation. The stochastic gradient descent method evaluates the gradient of the objective function only at a few samples for each iteration, which reduces the computational cost significantly. Various numerical experiments are conducted to illustrate the efficiency of the method and significant increases of the absorptance for the optimal random structures. We also examine the convergence of the stochastic gradient descent algorithm theoretically and prove that the numerical method is convergent under certain assumptions for the random interfaces.


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Key words: Optimal design, random rough surface, solar cell, Helmholtz equation, stochastic gradient descent method.

## 1 Introduction

Thin-film silicon solar cell is an attractive photovoltaic device because it attains a small thickness, which results in significant savings of material and energy during the fabrication. The cell consists of hydrogenated amorphous silicon (a-Si:H) as the absorbing layer,

[^0]| Glass |
| :---: |
| Transparent Conductive Oxide (TCO) |
| $\mathrm{a}-\mathrm{Si}: \mathrm{H}$ |
| TCO |
| Al |

Figure 1: A schematic plot of thin-film solar cells.
sandwiched between the transparent conductive oxide (TCO) layers for conducting the electric current. Figure 1 shows the structure of a typical thin-film solar cell, wherein the glass substrate on the top allows the incoming light to enter the cell and the highly reflective aluminum contact layer at the bottom enhances the absorption of light within the cell.

The a-Si:H layer in the thin-film solar cell is sufficiently absorptive at smaller optical wavelengths but poorly absorptive at larger wavelengths (typically $>600 \mathrm{~nm}$ ), which is responsible for the low overall efficiency of the cell. One way to increase the absorption within the solar cell and enhance its performance is to engineer the structure by texturing the interfaces between the different layers in a random manner $[1,10,11,13,18,22]$. The randomly textured surfaces lower the reflection losses at the entrance facet and scatter the light, thereby increasing the optical path of each photon in the solar cell. In realistic fabrication, the surfaces of the TCO layers in Figure 1 are textured randomly, which is achieved at low cost by controlling the deposition parameter of TCO films sputtered on substrates [16]. We would also like to point out several other ways to increase the absorption efficiency of solar cells, such as anti-reflection coating, dielectric gratings, and plasmonic nanoparticles [ $3,6,9,12,19$ ], although these techniques may be costly in fabrication.

The design and optimization of random surfaces in thin-film solar cells are mostly performed by the ad hoc procedures, where one computes the absorptance of the cell for chosen statistical parameters and obtains the optimal parameters from the comparison of the computed absorptance values $[10,11,13,18]$. Such ad hoc schemes are computationally inefficient and the optimal solutions heavily depend on the set of statistical parameters being chosen. To provide a systematic computational framework, in [4] we formulate the optimal design of random surface textures as a random PDE-constrained problem and apply the gradient-based algorithm to solve the optimization problem. The optimization problem seeks to maximize the mean absorptance function for the solar cells by sampling random surfaces in the appropriate probability space. We employ the Monte-Carlo method for sampling the probability space in [4] and apply the adjoint state method for computing the gradient at each sample. The optimal random textures give
rises to significant absorption enhancement, with the photon absorptance much higher than the existing random textures.

Albeit being able to provide optimal random surface textures, the numerical algorithm based on the Monte-Carlo sampling and the adjoint state method is computationally expensive, due to the largeness of the samples needed in computing the gradient average and the necessity to solve the underlying governing PDEs to obtain the gradient for each sample. In this work, we adopt the stochastic gradient descent method, which is a key ingredient of machine learning algorithms (cf. [7]), to solve the stochastic optimization problems. The new algorithm can obtain the statistical parameters of the optimal random textures, and its computational cost is significantly lower compared to the full gradient descent approach. We show that the numerical method is convergent under certain assumptions on the step sizes of the iterative algorithm and the random interfaces. In addition, in contrast to the optimization of one single random interface in [4], we consider the optimization of several random interfaces as well as the optimization of the random boundary for the solar cell. We also investigate more sophisticated configurations when optimization is performed over a frequency band or with multiple incident angles, which are computationally formidable if one attempts to solve by the full gradient method developed in [4].


Figure 2: Schematic plot of the multi-layered medium in the reference periodic cell with $0<x_{1}<\Lambda$. The layers from the bottom to the top are $D_{1}, D_{2}, \cdots, D_{\ell}$. The boundary at the bottom is given by $\Gamma_{1}$ and the interface between the two layers $D_{j-1}$ and $D_{j}$ is given by $\Gamma_{j}(j=2, \cdots, \ell)$.

More specifically, we consider the multi-layered structure in $\mathbb{R}^{2}$ as depicted in Figure 2 , which consists of several layers $D_{1}, D_{2}, \cdots, D_{\ell}$ from the bottom to the top. The boundary at the bottom of the structure $\Gamma_{1}$ and the interface $\Gamma_{j}(j=2, \cdots, \ell)$ between the two layers $D_{j-1}$ and $D_{j}$ are textured randomly. For each random sample $\zeta$, the interface $\Gamma_{j}$ is represented as $\Gamma_{j}(\zeta):=\left\{\left(x_{1}, x_{2}\right) \mid x_{2}=f_{j}\left(\zeta, \alpha_{j} ; x_{1}\right)\right\}$, wherein $\alpha_{j} \in \mathbb{R}^{d}$ represents the sta-
tistical parameters of the interface and $f_{j}$ is the profile function for the interface. The optimization problem is to solve for the optimal statistical parameters $\left\{\alpha_{j}\right\}_{j=1}^{\ell}$ such that the overall absorbtance within these layers is maximized. The problem setup will be discussed in more details in Section 2.

The rest of the paper is organized as follows. In Section 2 we introduce the mathematical model for the optical scattering problem by random rough surfaces and formulate the optimal design problems. The shape derivatives and the gradient of the objective function are derived by the adjoint state method in Section 3. We present the stochastic gradient descent method for the optimization problems and examine the convergence of the method in Section 4 . Finally, various numerical experiments are given in Section 5 to demonstrate the efficiency of the numerical method.

## 2 Mathematical formulation of the optimal design problems

### 2.1 Mathematical model for optical scattering problem by random surfaces

We assume that the whole structure is periodic along the $x_{1}$ direction with the period $\Lambda$, considering that the solar cells are usually arranged periodically in fabrication. For each random sample $\zeta$ and for $j=1,2, \cdots, \ell-1$, we let

$$
\begin{equation*}
D_{j}(\zeta):=\left\{\left(x_{1}, x_{2}\right): 0<x_{1}<\Lambda_{1}, f_{j}\left(\zeta ; x_{1}\right)<x_{2}<f_{j+1}\left(\zeta ; x_{1}\right)\right\} \tag{2.1}
\end{equation*}
$$

be the $j$ th layer in the reference period shown in Figure 2, and

$$
\begin{equation*}
D_{\ell}(\zeta):=\left\{\left(x_{1}, x_{2}\right): 0<x_{1}<\Lambda_{1}, x_{2}>f_{\ell}\left(\zeta ; x_{1}\right)\right\} \tag{2.2}
\end{equation*}
$$

be the domain on the top. For each $j$, the interface profile function satisfies

$$
f_{j}\left(\zeta ; x_{1}+\Lambda\right)=f_{j}\left(\zeta, x_{1}\right) \quad \text { for }-\infty<x_{1}<\infty,
$$

and $f_{j}\left(\zeta ; x_{1}\right)$ is a stationary random process in the reference period with $0<x_{1}<\Lambda$. This will be elaborated in Section 2.2 .

The relative permittivity function $\varepsilon_{r}$ attains the value $\varepsilon_{r, j}$ in each layer $D_{j}$. We consider the transverse electric (TE) polarization for the optical wave, in which the electric field attains the form $E=(0,0, u)$. The structure is illuminated by a time-harmonic incident plane wave $u^{i}=e^{i k_{0} q_{\ell}(\sin \theta,-\cos \theta) \cdot x}$, where $k_{0}$ be the free-space wavenumber, $\theta \in\left(-\frac{\pi}{2}, \frac{\pi}{2}\right)$ is the incident angle and $q_{\ell}:=\sqrt{\varepsilon_{r, \ell}}$ represents the refractive index in $D_{\ell}$. For simplicity of notation, here and henceforth, we set the wavenumber in $D_{\ell}$ as $k_{\ell}=k_{0} q_{\ell}$ and express the incident wave as $u^{i}=e^{i\left(\tau x_{1}-\rho x_{2}\right)}$, in which $\tau=k_{\ell} \sin \theta$ and $\rho=k_{\ell} \cos \theta$ is the wavenumber in the horizontal and vertical direction respectively. The total field $u$ after the scattering consists of the incident wave $u^{i}$ and the diffracted wave $u^{s}$. For each sample $\zeta \in \Omega$, the total field $u$ satisfies

$$
\begin{equation*}
\Delta u(\zeta ; x)+k_{0}^{2} \varepsilon_{r, j} u(\zeta ; x)=0 \quad \text { for } x \in D_{j}(\zeta), j=1,2, \cdots, \ell . \tag{2.3}
\end{equation*}
$$

Along the interfaces $\Gamma_{j}(\zeta)=\left\{\left(x_{1}, x_{2}\right) \mid 0<x_{1}<\Lambda_{1}, x_{2}=f_{j}\left(\zeta ; x_{1}\right)\right\}$ for $j=2, \cdots, \ell$, there hold

$$
\begin{align*}
u_{+}\left(\zeta ; x_{1}, f_{j}\left(\zeta, x_{1}\right)\right) & =u_{-}\left(\zeta ; x_{1}, f_{j}\left(\zeta, x_{1}\right)\right),  \tag{2.4}\\
\partial_{v} u_{+}\left(\zeta ; x_{1}, f_{j}\left(\zeta, x_{1}\right)\right) & =\partial_{v} u_{-}\left(\zeta ; x_{1}, f_{j}\left(\zeta, x_{1}\right)\right), \tag{2.5}
\end{align*}
$$

which follow from the continuity of the electric field and magnetic field. In the above, $v$ denotes the unit normal vector along $\Gamma_{j}$ pointing toward $D_{j}, u_{ \pm}$and $\partial_{\nu} u_{ \pm}$denote the limits of $u$ and $\partial_{v} u$ from above and below the surface, respectively. In addition, due to periodicity of the medium along the $x_{1}$ direction, we impose the quasi-boundary condition on the boundary walls of the periodic cell (cf. [2]):

$$
\begin{equation*}
u\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} u\left(\zeta ; 0, x_{2}\right) \quad \text { for } x_{2}>f_{1}(\zeta ; 0), \tag{2.6}
\end{equation*}
$$

in which $\tau$ is the horizontal wavenumber defined above. For a perfectly conducting contact layer $D_{1}$ such as aluminum depicted in Figure 1), there holds

$$
\begin{equation*}
u\left(\zeta ; x_{1}, f_{1}\left(\zeta ; x_{1}\right)\right)=0, \quad 0<x_{1}<\Lambda, \tag{2.7}
\end{equation*}
$$

along the boundary $\Gamma_{1}$. This implies that the optical light is totally reflected to the cell and no light is transmitted through $\Gamma_{1}$.

By virtue of the quasi-periodicity boundary condition, the solution to (2.3) - (2.7) can be expressed as a sum of a Fourier series. In particular, in the domain $D_{\ell}$, the diffracted field $u^{s}$ attains the so-called Rayleigh expansion (cf. [2,8]):

$$
\begin{equation*}
u^{s}\left(\zeta ; x_{1}, x_{2}\right)=\sum_{n=-\infty}^{\infty} \hat{u}_{n}^{s}(\zeta ; b) e^{i \kappa_{n} x_{1}+i \eta_{n}\left(x_{2}-b\right)} \quad \text { for } x_{2} \geq b \tag{2.8}
\end{equation*}
$$

where $b>\max _{-\infty<x_{1}<\infty} f_{\ell}\left(x_{1}\right)$ is a constant, $\kappa_{n}:=\tau+\frac{2 \pi n}{\Lambda}$ for $n \in \mathbb{Z}$, and

$$
\eta_{n}= \begin{cases}\sqrt{k_{\ell}^{2}-\kappa_{n}^{2}}, & k_{\ell}>\kappa_{n}  \tag{2.9}\\ i \sqrt{\kappa_{n}^{2}-k_{\ell}^{2}}, & k_{\ell}<\kappa_{n}\end{cases}
$$

The Fourier mode $e^{i \kappa_{n} x_{1}+i \eta_{n}\left(x_{2}-b\right)}$ is called the $n$th diffraction order and the corresponding Fourier coefficient $\hat{u}_{n}^{s}(\zeta ; b)$ is defined by

$$
\begin{equation*}
\hat{u}_{n}^{s}(\zeta ; b)=\frac{1}{\Lambda} \int_{0}^{\Lambda} u^{s}\left(\zeta ; x_{1}, b\right) e^{-i k_{n} x_{1}} d x_{1} . \tag{2.10}
\end{equation*}
$$

Here we assume that $\kappa_{n} \neq k_{\ell}$ to exclude resonances. Then we can introduce the Dirichlet-to-Neumann map $T$ on the line $x_{2}=b$ as

$$
\begin{equation*}
\frac{\partial u^{s}}{\partial x_{2}}\left(\zeta ; x_{1}, b\right)=\sum_{n=-\infty}^{\infty} i \eta_{n} \hat{u}_{n}^{s}(\zeta ; b) e^{i i_{n} x_{1}}=: T\left[u^{s}\left(\zeta ; x_{1}, b\right)\right] . \tag{2.11}
\end{equation*}
$$

Since $u=u^{i}+u^{s}$, there holds

$$
\begin{equation*}
\frac{\partial u}{\partial x_{2}}\left(\zeta ; x_{1}, b\right)=T\left(u\left(\zeta ; x_{1}, b\right)\right)+g \tag{2.12}
\end{equation*}
$$

where $g=-2 i \rho e^{i \tau x_{1}-i \rho b}$.
Let $\Gamma(\zeta)=\bigcup_{j=2}^{\ell} \Gamma_{j}(\zeta)$ and

$$
D(\zeta):=\left\{\left(x_{1}, x_{2}\right): 0<x_{1}<\Lambda_{1}, f_{1}\left(\zeta ; x_{1}\right)<x_{2}<b\right\}
$$

In light of (2.3) - (2.7) and (2.12), for each sample $\zeta$, the total field $u$ satisfies the following boundary value problem in the domain $D$ :

$$
\begin{cases}\Delta u(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} u(\zeta ; \cdot)=0 & \text { in } D(\zeta) \backslash \Gamma(\zeta)  \tag{2.13}\\ u\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} u\left(\zeta ; 0, x_{2}\right), & f_{1}(\zeta ; 0)<x_{2}<b \\ u\left(\zeta ; x_{1}, f_{1}\left(\zeta ; x_{1}\right)\right)=0, & 0<x_{1}<\Lambda \\ \frac{\partial u}{\partial x_{2}}\left(\zeta ; x_{1}, b\right)=T\left(u\left(\zeta ; x_{1}, b\right)\right)+g, & 0<x_{1}<\Lambda\end{cases}
$$

In addition, $u$ satisfies the conditions (2.4) - (2.5) along the interfaces.

### 2.2 Representation of random surfaces

For each random interface $\Gamma_{j}$, we assume that its profile function $f_{j}=f_{j}\left(\zeta ; x_{1}\right)$ is a stationary random process for $x_{1} \in[0, \Lambda]$, with a continuous and bounded covariance function $C_{j}\left(x_{1}, \tilde{x}_{1}\right)=c_{j}\left(x_{1}-\tilde{x}_{1}\right)$. We consider the Gaussian type covariance function with

$$
c_{j}\left(x_{1}-\tilde{x}_{1}\right)=\left(\alpha_{j}^{(1)}\right)^{2} \exp \left(-\left|x_{1}-\tilde{x}_{1}\right|^{2} /\left(\alpha_{j}^{(2)}\right)^{2}\right)
$$

where $\alpha_{j}^{(1)}$ is the root mean square and $\alpha_{j}^{(2)}$ is the correlation length of the surface $\Gamma_{j}(\zeta)$ satisfying $0<\alpha_{j}^{(2)} \ll \Lambda$. Such a covariance function is usually used for the modeling of rough surfaces [17].

By the Karhunen-Loève expansion (cf. [14]), the random process $f_{j}\left(\zeta ; x_{1}\right)$ can be represented as

$$
f_{j}\left(\zeta ; x_{1}\right)=f_{j}^{a}+\sum_{p=1}^{\infty} \sqrt{\lambda_{j p}} \xi_{j p} \varphi_{j p}\left(x_{1}\right)
$$

where $f_{j}^{a}$ is the average height of $f_{j}, \xi_{j p}$ are mutually uncorrelated random variables with zero mean and unit covariance, $\lambda_{j p}$ and $\varphi_{j p}(p=1,2, \cdots$,$) are the eigenvalues and$ eigenfunctions of covariance operator

$$
[K \varphi]\left(x_{1}\right):=\int_{0}^{\Lambda} c_{j}\left(x_{1}-\tilde{x}_{1}\right) \varphi\left(\tilde{x}_{1}\right) d \tilde{x}_{1}
$$

Since the covariance function $c_{j}\left(x_{1}\right)$ is even, we expand it as

$$
c_{j}\left(x_{1}\right)=\left(\alpha_{j}^{(1)}\right)^{2}\left[\frac{\widehat{c}_{j 0}}{2}+\sum_{p=1}^{\infty} \widehat{c}_{j p} \cos \left(\frac{2 p \pi x_{1}}{\Lambda}\right)\right] \quad \text { for } x_{1} \in[0, \Lambda]
$$

where $\widehat{c}_{j 0}, \widehat{c}_{j 1}, \widehat{c}_{j 2}, \cdots$, are the Fourier cosine expansion coefficients of the function $\exp \left(-x_{1}^{2} /\left(\alpha_{j}^{(2)}\right)^{2}\right)$. It can be shown that the covariance operator attains the eigenvalues

$$
\lambda_{j p}=\frac{\left(\alpha_{j}^{(1)}\right)^{2} \Lambda \hat{c}_{j p}}{2}, \quad p=0,1,2, \cdots
$$

The corresponding eigenfunctions are

$$
\varphi_{j p}\left(x_{1}\right)= \begin{cases}\sqrt{\frac{1}{\Lambda}}, & p=0 \\ \sqrt{\frac{2}{\Lambda}} \cos \left(\frac{2 p \pi x_{1}}{\Lambda}\right), & p>1 \text { and even } \\ \sqrt{\frac{2}{\Lambda}} \sin \left(\frac{2 p \pi x_{1}}{\Lambda}\right), & p>1 \text { and odd }\end{cases}
$$

for all $j$. Hence the Karhunen-Loève representation of the random process $f_{j}\left(\zeta ; x_{1}\right)$ is given by

$$
\begin{align*}
f_{j}\left(\zeta ; \alpha_{j} ; x_{1}\right)= & f_{j}^{a}+\sqrt{\lambda_{j 0}} \xi_{0}(\zeta) \sqrt{\frac{1}{\Lambda}} \\
& +\sum_{p=1}^{\infty} \sqrt{\lambda_{j p}}\left[\xi_{p}^{s}(\zeta) \sqrt{\frac{2}{\Lambda}} \sin \left(\frac{2 p \pi x_{1}}{\Lambda}\right)+\xi_{p}^{c}(\zeta) \sqrt{\frac{2}{\Lambda}} \cos \left(\frac{2 p \pi x_{1}}{\Lambda}\right)\right] \tag{2.14}
\end{align*}
$$

where $\xi_{0}, \xi_{p}^{s}$ and $\xi_{p}^{c}$ are mutually uncorrelated random variables with zero mean and unit covariance. $\alpha_{j}=\left(\alpha_{j}^{(1)}, \alpha_{j}^{(2)}\right)$ represents the statistical parameters of the interface. We express the explicit dependence of $f_{j}$ on $\alpha_{j}$ here and afterwards when necessary. In particular, when $f_{j}=f_{j}\left(\zeta ; x_{1}\right)$ is a stationary Gaussian process, $\xi_{0}, \xi_{p}^{s}$ and $\xi_{p}^{c}$ are independent and identically distributed Gaussian random variables with zero mean and unit covariance.

A finite-term Karhunen-Loève expansion is used in the computation so that the remaining terms are sufficiently small. Since the eigenvalues $\left\{\lambda_{j p}\right\}_{j=0}^{\infty}$ converge to 0 fast for the given smooth kernel $c\left(x_{1}-\tilde{x}_{1}\right)$, such an approximation yields high-order accuracy with a small number of terms in the expansion. Therefore, here and henceforth, for simplicity we use a finite-term approximation of (2.14) with $p \leq P_{0}$. By letting $\lambda_{j p}=\left(\alpha_{j}^{(1)}\right)^{2} \bar{\lambda}_{j p}$
with $\bar{\lambda}_{j p}=\frac{\Lambda \hat{c}_{j p}}{2}$, we express the profile of the random surface by

$$
\begin{equation*}
f_{j}\left(\zeta ; \alpha_{j} ; x_{1}\right)=f_{j}^{a}+\alpha_{j}^{(1)} \cdot \bar{f}_{j}\left(\zeta ; \alpha_{j}^{(2)} ; x_{1}\right) \tag{2.15}
\end{equation*}
$$

where

$$
\begin{align*}
\bar{f}_{j}\left(\zeta ; \alpha_{j}^{(2)} ; x_{1}\right)= & \sqrt{\bar{\lambda}_{j 0}} \xi_{0}(\zeta) \sqrt{\frac{1}{\Lambda}}+\sum_{p=1}^{P_{0}} \sqrt{\bar{\lambda}_{j p}}\left[\xi_{p}^{s}(\zeta) \sqrt{\frac{2}{\Lambda}} \sin \left(\frac{2 p \pi x_{1}}{\Lambda}\right)\right. \\
& \left.+\xi_{p}^{c}(\zeta) \sqrt{\frac{2}{\Lambda}} \cos \left(\frac{2 p \pi x_{1}}{\Lambda}\right)\right] \tag{2.16}
\end{align*}
$$

and it is independent of the root mean square $\alpha_{j}^{(1)}$. For each random sample, $f_{j}\left(\zeta ; \alpha_{j}, \cdot\right)$ is a smooth function and depends continuously on the statistical parameters $\alpha_{j}$.

### 2.3 Optimal design problems

For each sample $\zeta \in \Omega$, in light of the Rayleigh expansion (2.8), the diffracted field can be rewritten as

$$
u^{s}(\zeta ; \cdot)=\sum_{n=-\infty}^{\infty} r_{n}(\zeta) e^{i k_{n} x_{1}+i \eta_{n} x_{2}}
$$

where the reflection coefficient $r_{n}(\zeta)=\hat{u}_{n}^{s}(\zeta ; b) e^{-i \eta_{n} b}$, and $\hat{u}_{n}^{s}$ are the Fourier coefficients of the diffracted field $u^{s}$ as defined in (2.10). Since $u=u^{i}+u^{s}, r_{n}(\zeta)$ can also be written as

$$
r_{n}(\zeta)= \begin{cases}\hat{u}_{n}(\zeta ; b) e^{-i \eta_{n} b}, & n \neq 0  \tag{2.17}\\ \hat{u}_{n}(\zeta ; b) e^{-i \rho b}-e^{-2 i k_{\ell} b}, & n=0\end{cases}
$$

where $\hat{u}_{n}(\zeta ; b)$ are the Fourier coefficients of the total field $u(\zeta ; \cdot)$ on $x_{2}=b$.
Let $\mathcal{N}:=\left\{n \in \mathbb{Z} \mid k_{\ell}^{2}-\kappa_{n}^{2}>0\right\}$ be the set of indices for all propagating modes in the Rayleigh expansion. The goal of optimal design is to trap the energy in the layers $D_{1}, \cdots, D_{\ell-1}$ as much as possible. In other words, we aim to minimize the energy that is being reflected to $D_{\ell}$. Let $\boldsymbol{\alpha}=\left(\alpha_{1}, \cdots, \alpha_{\ell}\right)^{\top}$ be the design variables, where $\alpha_{j}=\left(\alpha_{j}^{(1)}, \alpha_{j}^{(2)}\right)$ are the statistical parameters of the interface $\Gamma_{j}$ for $j=1, \cdots, \ell$. Using the reflection coefficients above, the reflectivity associated with the optical structure for each sample $\zeta$ is defined by

$$
R(\zeta ; \boldsymbol{\alpha})=\sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}}\left|r_{n}(\zeta)\right|^{2}
$$

where $\eta_{n}$ is defined in (2.9). The mean reflectivity is

$$
\begin{equation*}
E[R(\zeta ; \alpha)]:=\int_{\Omega} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}}\left|r_{n}(\zeta)\right|^{2} d P(\zeta), \tag{2.18}
\end{equation*}
$$

in which $\Omega$ and $P$ denotes the random sample space and the probability measure, respectively.

Let $Q(\boldsymbol{\alpha}):=E[R(\zeta ; \boldsymbol{\alpha})]$, the optimal design problem for the fixed wavenumber $k_{0}$ is to minimize the mean reflectivity $Q(\boldsymbol{\alpha})$ by solving the following stochastic optimization problem over an admissible set $U_{\alpha}:=\left(\left(0, \alpha_{1}^{\max }\right) \times\left(0, \alpha_{2}^{\max }\right)\right)^{\ell}$ :

Problem (I)

$$
\begin{equation*}
\min _{\alpha \in U_{\alpha}} Q(\boldsymbol{\alpha}) \tag{2.19}
\end{equation*}
$$

Since the solar frequency spectrum ranges from 300 nm to about 3000 nm , and the angle of the incidence for the incoming light changes during the daytime, it is also important to investigate the corresponding optimization problems in these realistic scenarios. In the case of optimal design over a frequency band, assuming that the wavelength $\lambda$ for the incident wave is within the range $\left[\lambda_{\min }, \lambda_{\max }\right]$, the corresponding stochastic optimization problem is cast as

Problem (II) $\min _{\boldsymbol{\alpha} \in U_{\boldsymbol{\alpha}}} Q(\boldsymbol{\alpha})$, where $Q(\boldsymbol{\alpha}):=E\left[\int_{\lambda_{\min }}^{\lambda_{\max }} R(\zeta ; \boldsymbol{\alpha}, \lambda) d \lambda\right]$.
Note that in this configuration, the refractive index $\varepsilon_{r}(\lambda ; x)$ is a function of the wavelength $\lambda$. Finally, the optimal design problem with multiple incident angles with $\theta \in\left[\theta_{\min }, \theta_{\max }\right]$ is formulated as follows:

$$
\begin{equation*}
\text { Problem (III) } \quad \min _{\boldsymbol{\alpha} \in U_{\alpha}} Q(\boldsymbol{\alpha}), \text { where } Q(\boldsymbol{\alpha}):=E\left[\int_{\theta_{\min }}^{\theta_{\max }} R(\zeta ; \boldsymbol{\alpha}, \theta) d \theta\right] \tag{2.21}
\end{equation*}
$$

Both Problem (II) and (III) are computationally more expensive than Problem (I) due to the necessity to sample over the frequency band or at different incident angles. It is computationally formidable by using the gradient descent algorithm in [4] directly.

## 3 The computation of the gradient $D_{\alpha} R(\zeta ; \alpha)$

To perform the optimization, one needs to compute the gradient of the objective function. In this section, we derive the gradient $D_{\alpha} R(\zeta ; \boldsymbol{\alpha})$ of the reflectivity $R(\zeta ; \boldsymbol{\alpha})$ at each sample. The shape derivatives are obtained by analyzing the sensitivity of the reflectivity $R$ upon the perturbation of the interface/surface. Note that $\Gamma_{1}$ at the bottom is the boundary of the structure, while $\Gamma_{2}, \cdots, \Gamma_{\ell}$ are interfaces between two layers. Thus the derivations of the shape derivatives are different for the boundary $\Gamma_{1}$ and the rest of interfaces. We present the shape derivative formulas for $\Gamma_{1}$ and $\Gamma_{j}(j \geq 2)$ respectively in the following theorem and give the detailed proof in Section 3.2.1 and 3.2.2. The readers are referred to [21] for analysis of derivatives in various shape optimizations.

Theorem 3.1. Denote $D_{\alpha_{j}}:=\left(\frac{\partial}{\partial \alpha_{j}^{(1)}}, \frac{\partial}{\partial \alpha_{j}^{(2)}}\right)$ for $j=1,2, \cdots, \ell$. For each sample $\zeta$,

$$
\begin{equation*}
D_{\alpha_{1}} R(\zeta ; \boldsymbol{\alpha})=\frac{2}{\Lambda} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[\left(\hat{u}_{n}(\zeta ; b)-a_{n} e^{-2 i k_{l} b+i \rho b}\right) \cdot \int_{0}^{\Lambda}\left(\frac{\partial \bar{u}}{\partial v} \cdot \frac{\partial u_{n}^{*}}{\partial v} v_{2}\right) \cdot D_{\alpha_{1}} f_{1} d x_{1}\right] \tag{3.1a}
\end{equation*}
$$

and

$$
\begin{align*}
& D_{\alpha_{j}} R(\zeta ; \alpha)=\frac{2 k_{0}^{2}}{\Lambda} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[\left(\hat{u}_{n}(\zeta ; b)-a_{n} e^{-2 i k_{e} b+i \rho b}\right) \cdot \overline{\left(\varepsilon_{r, j}-\varepsilon_{r, j-1}\right)}\right. \\
&\left.\left.\cdot \int_{0}^{\Lambda}\left[\bar{u} u_{n}^{*}\right]\right|_{\left(x_{1}, f_{j}\right)} \cdot D_{\alpha_{j}} f_{j} d x_{1}\right], \quad j=2, \cdots, \ell . \tag{3.1b}
\end{align*}
$$

Here $a_{0}=1$ and $a_{n}=0$ if $n \neq 0$, and $v=\left(v_{1}, v_{2}\right)^{\top}$ is the unit normal vector pointing to the interior of $D$ along $\Gamma_{1}$. $u$ is the solution to the forward problem (2.13) and $u_{n}^{*}$ solves the following adjoint problem

$$
\left\{\begin{array}{l}
\Delta u_{n}^{*}(\zeta ; \cdot)+k_{0}^{2} \bar{\varepsilon}_{r} u_{n}^{*}(\zeta ; \cdot)=0 \text { in } D(\zeta) \backslash \Gamma(\zeta),  \tag{3.2}\\
u_{n}^{*}\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} u_{n}^{*}\left(\zeta ; 0, x_{2}\right), \quad f_{1}(\zeta ; 0)<x_{2}<b, \\
u_{n}^{*}\left(\zeta ; x_{1}, f_{1}\left(\zeta ; x_{1}\right)\right)=0, \quad 0<x_{1}<\Lambda, \\
\partial_{x_{2}} u_{n}^{*}\left(\zeta ; x_{1}, b\right)=T^{*}\left(u\left(\zeta ; x_{1}, b\right)\right)+e^{i k_{n} x_{1}} .
\end{array}\right.
$$

In the above theorem, $\left.\left[u_{n}^{*} \bar{u}\right]\right|_{\left(x_{1}, f_{j}\right)}$ denotes the restriction of $u_{n}^{*} \bar{u}$ to the surface $\Gamma_{j}(\zeta)$. $T^{*}$ is the adjoint operator of $T$ such that

$$
\langle T u, v\rangle=\left\langle u, T^{*} v\right\rangle,
$$

where $\langle\cdot, \cdot\rangle$ stands for the inner product over the function space $L^{2}(0, \Lambda)$.

### 3.1 Proof of formula (3.1a)

Let $H_{\tau}^{1}(D):=\left\{u \in H^{1}(D): u=0\right.$ on $\left.\Gamma_{1}, u\left(\Lambda, x_{2}\right)=e^{i \tau \Lambda} u\left(0, x_{2}\right)\right\}$, where the function space

$$
H^{1}(D)=\left\{u(x):\left(\int_{D}|u(x)|^{2} d x\right)^{\frac{1}{2}}<\infty,\left(\int_{D}|\nabla u(x)|^{2} d x\right)^{\frac{1}{2}}<\infty\right\} .
$$

We introduce the bilinear form

$$
a(u, w):=\int_{D} \nabla u \cdot \nabla \bar{w}-k_{0}^{2} \varepsilon_{r} u \bar{w} d x-\langle T u, w\rangle \quad \text { for } u, w \in H_{\tau}^{1}(D) .
$$

Here $\langle\cdot, \cdot\rangle$ stands for the inner product over the function space $L^{2}(0, \Lambda)$. Then each random sample $\zeta$, the weak solution $u(\zeta ; \cdot) \in H_{\tau}^{1}(D)$ for the boundary value problem (2.13) satisfies

$$
\begin{equation*}
a(u(\zeta ; \cdot), w)=\langle g, w\rangle \quad \text { for all } w \in H_{\tau}^{1}(D) \tag{3.3}
\end{equation*}
$$

For each random sample, define the mapping $S: \alpha_{1}=\left(\alpha_{1}^{(1)}, \alpha_{1}^{(2)}\right) \rightarrow u\left(\zeta ; x_{1}, b\right)$, where $u$ is the solution to boundary value problem (2.13). We let $D_{\alpha_{1}} S:=\left(\frac{\partial S}{\partial \alpha_{1}^{(1)}}, \frac{\partial S}{\partial \alpha_{1}^{(2)}}\right)$, in which

$$
\frac{\partial S}{\partial \alpha_{1}^{(j)}}=\lim _{\delta \rightarrow 0} \frac{S\left(\alpha^{(1)}+\delta e_{j}\right)-S\left(\alpha^{(1)}\right)}{\delta} \quad \text { for } j=1,2
$$

and $e_{j}$ is the unit vector.
Lemma 3.2. The derivative $D_{\alpha_{1}} S$ exits and $\frac{\partial S}{\partial \alpha_{1}^{(j)}}=u_{j 0}\left(\zeta ; x_{1}, b\right)$ for $j=1,2$, where $u_{j 0}$ solves

$$
\left\{\begin{array}{l}
\Delta u_{j 0}(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} u_{j 0}(\zeta ; \cdot)=0 \text { in } D(\zeta) \backslash \Gamma(\zeta),  \tag{3.4}\\
u_{j 0}\left(\zeta ; 0, x_{2}\right)=e^{i \tau \Lambda} u_{j 0}\left(\zeta ; \Lambda, x_{2}\right), \\
u_{j 0}\left(\zeta ; x_{1}, f_{1}\left(x_{1}\right)\right)=-\frac{\partial f_{1}}{\partial \alpha_{1}^{(j)}} \frac{\partial u}{\partial \nu} v_{2} ; 0<x_{1}<\Lambda \\
\partial_{x_{2}} u_{j 0}\left(\zeta ; x_{1}, b\right)=T\left(u_{j 0}\left(\zeta ; x_{1}, b\right)\right),
\end{array}\right.
$$

$v=\left(v_{1}, v_{2}\right)^{\top}$ is the unit normal vector pointing to $D$ along $\Gamma_{1}$, and $u$ is the solution to boundary value problem (2.13).

Proof. We only provide the proof for $\frac{\partial S}{\partial \alpha_{1}^{(1)}}$, and the proof for $\frac{\partial S}{\partial \alpha_{1}^{(2)}}$ is similar. Let $\alpha_{1}^{(1)}$ be perturbed by a small number $\delta$, then the new root mean square is $\left(\alpha_{1}^{(1)}\right)^{\delta}:=\alpha_{1}^{(1)}+\delta$ and the new boundary becomes $\Gamma_{1}^{\delta}:=\left\{\left(x_{1}, x_{2}\right): x_{2}=f_{1}^{\delta}\left(\zeta ; x_{1}\right)\right\}$, where $f_{1}^{\delta}=f_{1}+\delta \cdot \frac{\partial f_{1}}{\partial x_{1}^{(1)}}+$ $O\left(\delta^{2}\right)$. For simplicity of the notation, we introduce the vector function $W(x) \in C_{0}^{2}\left(\Gamma_{1} ; \mathbb{R}^{2}\right)$ such that $\delta \cdot W(x):=\left[0, f_{1}^{\delta}\left(x_{1}\right)-f_{1}\left(x_{1}\right)\right]^{\top}$ for $x \in \Gamma_{1}$, then the perturbed boundary can be expressed as $\Gamma_{1}^{\delta}=\left\{x+\delta \cdot W(x): x \in \Gamma_{1}\right\}$.

We denote the domain after the perturbation by $D^{\delta}$. The perturbed total field $u^{\delta}$ satisfies

$$
\left\{\begin{array}{l}
\Delta u^{\delta}(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} u^{\delta}(\zeta ; \cdot)=0 \text { in } D^{\delta}(\zeta) \backslash \Gamma(\zeta)  \tag{3.5}\\
u^{\delta}\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} u^{\delta}\left(\zeta ; 0, x_{2}\right) \\
u^{\delta}\left(\zeta ; x_{1}, f_{1}\left(\zeta ; x_{1}\right)\right)=0 ; 0<x_{1}<\Lambda \\
\partial_{x_{2}} u^{\delta}\left(\zeta ; x_{1}, b\right)=T\left(u^{\delta}\left(\zeta ; x_{1}, b\right)\right)+g
\end{array}\right.
$$

The weak solution $u^{\delta}$ for the above boundary value problem satisfies

$$
\begin{equation*}
a^{\delta}\left(u^{\delta}(\zeta ; \cdot), w^{\delta}\right)=\left\langle g, w^{\delta}\right\rangle \tag{3.6}
\end{equation*}
$$

for all $w^{\delta} \in H_{0}^{1}\left(D^{\delta}\right)$, where

$$
a^{\delta}\left(u^{\delta}(\zeta ; \cdot), w^{\delta}\right)=\int_{D^{\delta}} \nabla u^{\delta}(\zeta ; \cdot) \cdot \nabla \overline{w^{\delta}}-k_{0}^{2} \varepsilon_{r} u^{\delta}(\zeta ; \cdot) \overline{w^{\delta}} d x-\left\langle T u^{\delta}(\zeta ; \cdot), w^{\delta}\right\rangle .
$$

Let us extend the definition of $W(x)$ to the closure of the whole domain $D$ such that $W \in C^{2}\left(\bar{D} ; \mathbb{R}^{2}\right)$ and $W(x)=0$ on the boundary $x_{2}=b$. Correspondingly, we introduce a map $\psi$ from $D$ to $D^{\delta}$ by letting $x=\psi(y)=y+W(y)$ for $y \in D$. The inverse map of $\psi$ is denoted as $\phi(x)$, which maps $D^{\delta}$ to $D$. Let $\tilde{u}^{\delta}(y)=u^{\delta}(\psi(y)), \tilde{w}^{\delta}=w^{\delta}(\psi(y))$, then $\tilde{u}^{\delta}$ and $\tilde{w}$ are defined on $D$. It is straightforward to show that $\frac{\partial u^{\delta}}{\partial x_{1}}=\sum_{m=1}^{2} \frac{\partial \tilde{u}^{\delta}}{\partial y_{m}} \frac{\partial \phi_{m}}{\partial x_{1}}$, where $\phi_{1}, \phi_{2}$ are the two components of the mapping $\phi$. By change of variables, we obtain

$$
a^{\delta}\left(u^{\delta}(\zeta ; \cdot), w^{\delta}\right)=\int_{D}\left[\sum_{m, n=1}^{2} b_{m n} \frac{\partial \tilde{u}^{\delta}(\zeta ; \cdot)}{\partial y_{m}} \frac{\partial \overline{\tilde{w}^{\delta}}}{\partial y_{n}}-k_{0}^{2} \varepsilon_{r} \tilde{u}^{\delta}(\zeta ; \cdot) \overline{\tilde{w}^{\delta}}\right] J d y-\left\langle T \tilde{u}^{\delta}(\zeta ; \cdot), \tilde{w}^{\delta}\right\rangle
$$

where $J=\operatorname{det} \nabla \psi, b_{m n}=\sum_{i=1}^{2} \frac{\partial \phi_{m}}{\partial x_{i}} \frac{\partial \phi_{n}}{\partial x_{i}}$. Define a new bilinear form

$$
\tilde{a}^{\delta}\left(\tilde{u}^{\delta}, w\right)=\int_{D}\left[\sum_{m, n=1}^{2} b_{m n} \frac{\partial \tilde{u}^{\delta}}{\partial y_{m}} \frac{\partial \bar{w}}{\partial y_{n}}-k_{0}^{2} \varepsilon_{r} \tilde{u}^{\delta} \bar{w}\right] J d y-\left\langle T \tilde{u}^{\delta}, w\right\rangle
$$

for $\tilde{u}^{\delta}, w \in H_{\tau}^{1}(D)$. Then (3.6) is equivalent to finding $\tilde{u}^{\delta} \in H_{\tau}^{1}(D)$ such that

$$
\begin{equation*}
\tilde{a}\left(\tilde{u}^{\delta}(\zeta ; \cdot), w\right)=\langle g, w\rangle \tag{3.7}
\end{equation*}
$$

for all $w \in H_{\tau}^{1}(D)$.
From (3.3) and (3.7), it is seen that $\tilde{u}^{\delta}(\zeta ; \cdot)-u(\zeta ; \cdot)$ satisfies

$$
\begin{equation*}
a\left(\tilde{u}^{\delta}(\zeta ; \cdot)-u(\zeta ; \cdot), w\right)=-\left(\tilde{a}^{\delta}\left(\tilde{u}^{\delta}(\zeta ; \cdot), w\right)-a\left(\tilde{u}^{\delta}(\zeta ; \cdot), w\right)\right) \tag{3.8}
\end{equation*}
$$

For the right-hand side,

$$
\begin{align*}
\tilde{a}^{\delta}\left(\tilde{u}^{\delta}(\zeta ; \cdot), w\right)-a\left(\tilde{u}^{\delta}(\zeta ; \cdot), w\right)= & \int_{D}\left[\sum_{m, n=1}^{2} b_{m n} \frac{\partial \tilde{u}^{\delta}(\zeta ; \cdot)}{\partial y_{m}} \frac{\partial \bar{w}}{\partial y_{n}}-k_{0}^{2} \varepsilon_{r} \tilde{u}^{\delta}(\zeta ; \cdot \bar{w}] J d y\right.  \tag{3.9}\\
& -\int_{D} \nabla \tilde{u}^{\delta}(\zeta ; \cdot) \cdot \nabla \bar{w}-k_{0}^{2} \varepsilon_{r} \tilde{u}^{\delta}(\zeta ; \cdot) \bar{w} d y .
\end{align*}
$$

Let $V(x)$ be the leading-order of the vector function $W(x)$ and it is independent of $\delta$. Then it can be calculated that the Jacobian $J=1+\delta \nabla \cdot V+O\left(\delta^{2}\right)$, and $\left(b_{m n}\right) J=I-\delta\left(\tilde{b}_{m n}\right)+O\left(\delta^{2}\right)$, where $I$ is the $2 \times 2$ identity matrix and

$$
\begin{equation*}
\tilde{b}_{m n}=\nabla V+(\nabla V)^{T}-(\nabla \cdot V) I . \tag{3.10}
\end{equation*}
$$

Therefore,

$$
\begin{equation*}
a\left(\tilde{u}^{\delta}(\zeta ; \cdot)-u(\zeta ; \cdot), w\right)=\delta \int_{D} \sum_{m, n=1}^{2} \tilde{b}_{m n} \frac{\partial \tilde{u}^{\delta}(\zeta ; \cdot)}{\partial y_{m}} \frac{\partial \bar{w}}{\partial y_{n}}+k_{0}^{2} \varepsilon_{r}(\nabla \cdot V) u(\zeta ; \cdot) \bar{w} d y+O\left(\delta^{2}\right) \tag{3.11}
\end{equation*}
$$

Denote $u^{\prime}(\zeta ; \cdot)=\lim _{\delta \rightarrow 0} \frac{\tilde{u}^{\delta}(\zeta ; \cdot)-u(\zeta ; \cdot)}{\delta}$. Then $u^{\prime}(\zeta ; \cdot)$ satisfies the following variational formulation:

$$
\begin{equation*}
a\left(u^{\prime}(\zeta ; \cdot), w\right)=\int_{D} \sum_{m, n=1}^{2} \tilde{b}_{m n} \frac{\partial \tilde{u}^{\delta}(\zeta ; \cdot)}{\partial y_{m}} \frac{\partial \bar{w}}{\partial y_{n}}+k_{0}^{2} \varepsilon_{r}(\nabla \cdot V) u(\zeta ; \cdot) \bar{w} d y . \tag{3.12}
\end{equation*}
$$

By the formula (3.10), we have

$$
\begin{aligned}
\sum_{m, n=1}^{2} \tilde{b}_{m n} \frac{\partial u(\zeta ; \cdot)}{\partial y_{m}} \frac{\partial \bar{w}}{\partial y_{n}}= & \nabla(V \cdot \nabla \bar{w}) \cdot \nabla u(\zeta ; \cdot)+\nabla(V \cdot \nabla u(\zeta ; \cdot)) \cdot \nabla \bar{w} \\
& -\nabla \cdot[(\nabla u(\zeta ; \cdot) \cdot \nabla \bar{w}) V] .
\end{aligned}
$$

By the Green's formula and the boundary condition $u\left(\zeta_{;} \cdot\right)=w=0$ on $\Gamma_{1}$, there holds

$$
\int_{\Gamma_{1}}(V \cdot \nabla \bar{w}) \frac{\partial u(\zeta ; \cdot)}{\partial v}-(\nabla u(\zeta ; \cdot) \cdot \nabla \bar{w})(V \cdot v) d s=0 .
$$

Therefore, (3.12) can be reduced to

$$
\begin{align*}
a\left(u^{\prime}(\zeta ; \cdot), w\right)= & \int_{D}-\nabla(V \cdot \nabla \bar{w}) \cdot \nabla u(\zeta ; \cdot)+\nabla(V \cdot \nabla u(\zeta ; \cdot)) \cdot \nabla \bar{w}+k_{0}^{2} \varepsilon_{r}(\nabla \cdot V) u(\zeta ; \cdot) \bar{w} d y \\
& +\int_{\Gamma_{1}}(V \cdot \nabla \bar{w}) \frac{\partial u(\zeta ; \cdot)}{\partial v}-(\nabla u(\zeta ; \cdot) \cdot \nabla \bar{w})(V \cdot v) d s \\
= & \int_{D} k_{0}^{2} \varepsilon_{r} u(\zeta ; \cdot)(V \cdot \nabla \bar{w})+\nabla(V \cdot \nabla u(\zeta ; \cdot)) \cdot \nabla \bar{w}+k_{0}^{2} \varepsilon_{r}(\nabla \cdot V) u(\zeta ; \cdot \bar{w} d y \\
= & \int_{D} \nabla(V \cdot \nabla u(\zeta ; \cdot)) \cdot \bar{w}-k_{0}^{2} \varepsilon_{r}(V \cdot \nabla u(\zeta ; \cdot)) \bar{w} d y+\int_{D} k_{0}^{2} \varepsilon_{r} \nabla \cdot(u(\zeta ; \cdot) \bar{w} V) d y . \tag{3.13}
\end{align*}
$$

Since $\int_{D} k_{0}^{2} \varepsilon_{r} \nabla \cdot(u(\zeta ; \cdot) \bar{w} V) d y=0$ by the divergence theorem, we obtain
$a\left(u^{\prime}(\zeta ; \cdot), w\right)=\int_{D} \nabla(V \cdot \nabla u(\zeta ; \cdot)) \cdot \nabla \bar{w}-k_{0}^{2} \varepsilon_{r}(V \cdot \nabla \bar{u}(\zeta ; \cdot)) \bar{w} d y$ for any $w \in H_{\tau}^{1}(D) \cap H^{2}(D)$.
This implies that $u^{\prime}$ is a weak solution of the following boundary value problem:

$$
\left\{\begin{array}{l}
\Delta u^{\prime}(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} u^{\prime}(\zeta ; \cdot)=\left(\Delta+k_{0}^{2} \varepsilon_{r}\right)(V-\nabla u) \text { in } D(\zeta) \backslash \Gamma(\zeta), \\
u^{\prime}\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} u^{\prime}\left(\zeta ; 0, x_{2}\right), \\
u^{\prime}\left(\zeta ; x_{1}, f_{1}\left(x_{1}\right)\right)=0 ; 0<x_{1}<\Lambda, \\
\partial_{x_{2}} u^{\prime}\left(\zeta ; x_{1}, b\right)=T\left(u^{\prime}\left(\zeta ; x_{1}, b\right)\right) .
\end{array}\right.
$$

Let $u_{10}=u^{\prime}-V \cdot \nabla u$, then $u_{10}=u^{\prime}$ on $x_{2}=b$, and $u_{10}$ satisfies

$$
\left\{\begin{array}{l}
\Delta u_{10}(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} u_{10}(\zeta ; \cdot)=0 \text { in } D(\zeta) \backslash \Gamma(\zeta) \\
u_{10}\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} u_{10}\left(\zeta ; 0, x_{2}\right) \\
u_{10}\left(\zeta ; x_{1}, f_{1}\left(x_{1}\right)\right)=\frac{\partial f_{1}}{\partial \alpha_{1}^{(1)}} \frac{\partial u}{\partial v} v_{2} ; 0<x_{1}<\Lambda \\
\partial_{x_{2}} u_{10}\left(\zeta ; x_{1}, b\right)=T\left(u_{10}\left(\zeta ; x_{1}, b\right)\right)
\end{array}\right.
$$

This completes the proof of Lemma 3.2.

Next, let us prove formula (3.1a). From the definition of the reflectivity $R(\zeta ; \boldsymbol{\alpha})=$ $\sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}}\left|r_{n}(\zeta)\right|^{2}$, we have

$$
\frac{\partial R(\zeta ; \boldsymbol{\alpha})}{\partial \alpha_{1}^{(1)}}=2 \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[r_{n} \frac{\partial \bar{r}_{n}}{\partial \alpha_{1}^{(1)}}\right]
$$

where

$$
r_{n}(\zeta)= \begin{cases}\hat{u}_{n}(\zeta ; b) e^{-i \eta_{n} b}, & n \neq 0 \\ \hat{u}_{n}(\zeta ; b) e^{-i \rho b}-e^{-2 i k_{0} q_{\ell} b}, & n=0\end{cases}
$$

From a direct calculation, it follows that

$$
\frac{\partial r_{n}(\zeta)}{\partial \alpha_{1}^{(1)}}=\frac{\partial \hat{u}_{n}(\zeta ; b)}{\partial \alpha_{1}^{(1)}} e^{-i \eta_{n} b}
$$

Applying Lemma 3.2, we obtain $\frac{\partial u(\zeta ; b)}{\partial \alpha_{1}^{(1)}}=\lim _{\delta \rightarrow 0} \frac{u^{\delta}(\zeta ; b)-u(\zeta ; b)}{\delta}=u_{01}(\zeta ; b)$. Thus

$$
r_{n} \frac{\partial \bar{r}_{n}}{\partial \alpha_{1}^{(1)}}= \begin{cases}\hat{u}_{n}(\zeta ; b) \cdot \frac{1}{\Lambda} \int_{0}^{\Lambda} e^{i \kappa_{n} x_{1}} \overline{u_{01}\left(\zeta ; x_{1}, b\right)} d x_{1}, & n \neq 0  \tag{3.14}\\ \left(\hat{u}_{n}(\zeta ; b)-e^{-2 i k_{0} q_{\ell} b+i \rho b}\right) \cdot \frac{1}{\Lambda} \int_{0}^{\Lambda} e^{i \kappa_{n} x_{1}} \overline{u_{01}\left(\zeta ; x_{1}, b\right)} d x_{1}, & n=0\end{cases}
$$

Multiplying the differential equation in the adjoint problem (3.2) by $\overline{u_{01}(\zeta ; \cdot)}$ and multiplying the complex conjugate of the differential equation in (3.4) by $u_{n}^{*}(\zeta ; \cdot)$, and integrating over the domain $D_{j}$ for $j=1, \cdots, \ell$, we have

$$
\int_{D_{j}}\left(\Delta u_{n}^{*}(\zeta ; \cdot)+k_{0}^{2} \overline{\varepsilon_{r}} u_{n}^{*}(\zeta ; \cdot)\right) \overline{u_{01}(\zeta ; \cdot)}-u_{n}^{*}(\zeta ; \cdot) \overline{\left(\Delta u_{01}(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} u_{01}(\zeta ; \cdot)\right)} d x=0 ; \quad j=1, \cdots, \ell
$$

Applying the Green's formula for the above equations and adding them together, we obtain

$$
\begin{aligned}
& -\int_{0}^{\Lambda}\left(\left[0, \frac{\partial f_{1}}{\partial \alpha_{1}^{(1)}}\right]^{\top} \cdot v\right)\left[\frac{\partial \overline{u(\zeta ; \cdot)}}{\partial v} \cdot \frac{\partial u_{n}^{*}(\zeta ; \cdot)}{\partial v}\right] d x_{1} \\
& +\int_{0}^{\Lambda} \overline{u_{01}(\zeta ; \cdot)}\left(T^{*}\left(u_{n}^{*}(\zeta ; \cdot)+e^{i \kappa_{n} x_{1}}\right)\right)-u_{n}^{*}(\zeta ; \cdot) T\left(u_{01}(\zeta ; \cdot)\right) d x_{1}=0,
\end{aligned}
$$

where we have used the boundary conditions in the boundary value problems (3.2) and (3.4).

Since $T^{*}$ is the adjoint operator of $T$, there holds

$$
\int_{0}^{\Lambda} e^{i \kappa_{n} x_{1}} \overline{u_{01}\left(\zeta ; x_{1}, b\right)} d x_{1}=\int_{0}^{\Lambda}\left(\frac{\partial \overline{u(\zeta ; \cdot)}}{\partial v} \cdot \frac{\partial u_{n}^{*}(\zeta ; \cdot)}{\partial v} v_{2}\right) \frac{\partial f_{1}}{\partial \alpha_{1}^{(1)}} d x_{1} .
$$

Substituting into (3.14) leads to

$$
\frac{\partial R}{\partial \alpha_{1}^{(1)}}=\frac{2}{\Lambda} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} R e\left[\left(\hat{u}_{n}(\zeta ; b)-a_{n} e^{-2 i k_{0} q_{c} b+i \rho b}\right) \cdot \int_{0}^{\Lambda}\left(\frac{\partial \overline{u(\zeta ; \cdot)}}{\partial v} \cdot \frac{\partial u_{n}^{*}(\zeta ; \cdot)}{\partial v} v_{2}\right) \frac{\partial f_{1}}{\partial \alpha_{1}^{(1)}} d x_{1}\right] .
$$

Therefore,

$$
D_{\alpha_{1}} R=\frac{2}{\Lambda} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[\left(\hat{u}_{n}(\zeta ; b)-a_{n} e^{-2 i k_{0} q_{e} b+i \rho b}\right) \cdot \int_{0}^{\Lambda}\left(\frac{\partial \overline{u(\zeta ; \cdot)}}{\partial v} \cdot \frac{\partial u_{n}^{*}(\zeta ; \cdot)}{\partial v} v_{2}\right) \cdot D_{\alpha_{1}} f_{1} d x_{1}\right] .
$$

### 3.2 Proof of formula (3.1b)

To prove (3.1b), we need to derive the perturbation of the reflectivity $\delta R$ due to the perturbation of the interface by $\delta f_{j}$ induced by a small perturbation of $\alpha_{j}^{(1)}$ or $\alpha_{j}^{(2)}$ for $j=2, \cdots, \ell$. When the interface $\Gamma_{j}$ is perturbed as $f_{j}^{\delta}:=f_{j}+\delta f_{j}$, the permittivity $\varepsilon$ in $D_{j}$ becomes $\varepsilon_{r}^{\delta}:=\varepsilon_{r}+\delta \varepsilon_{r}$. It is observed that for any test function $v \in L^{2}(D)$, the inner product

$$
\left(v, \delta \varepsilon_{r}\right):=\int_{D} v(x) \overline{\delta \varepsilon_{r}}(x) d x=\int_{\operatorname{symdiff}\left(D_{j}, D_{j}^{\delta}\right)} v(x) \overline{\delta \varepsilon_{r}}(x) d x .
$$

Here $D_{j}$ and $D_{j}^{\delta}$ are the layers with the interfaces $f_{j}$ and $f_{j}^{\delta}$, respectively, and the symmetric difference of the two sets $D_{j}$ and $D_{j}^{\delta}$ is given by

$$
\operatorname{symdiff}\left(D_{j}, D_{j}^{\delta}\right)=\left(D_{j} \cup D_{j}^{\delta}\right) \backslash\left(D_{j} \cap D_{j}^{\delta}\right)
$$

Since the relative permittivity of the domain $D_{j-1}$ and $D_{j}$ are $\varepsilon_{r, j-1}$ and $\varepsilon_{r, j}$, respectively, the above inner product can be simplified as

$$
\begin{equation*}
\left(v, \delta \varepsilon_{r}\right)=\int_{0}^{\Lambda} v\left(x_{1}, f_{j}\left(x_{1}\right)\right) \overline{\left(\varepsilon_{r, j-1}-\varepsilon_{r, j}\right)} \cdot \delta f_{j} d x \tag{3.15}
\end{equation*}
$$

for an infinitesimal $\delta f$.
Let $\delta u$ denote the perturbation of the total field. As a result of perturbation analysis, $\delta u$ satisfies the following equations:

$$
\begin{cases}\Delta \delta u(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} \delta u(\zeta ; \cdot)=-k_{0}^{2} \delta \varepsilon_{r} u(\zeta ; \cdot) & \text { in } D(\zeta) \backslash \Gamma(\zeta),  \tag{3.16}\\ \delta u\left(\zeta ; \Lambda, x_{2}\right)=e^{i \tau \Lambda} \delta u\left(\zeta ; 0, x_{2}\right), & 0<x_{2}<b, \\ \delta u\left(\zeta ; x_{1}, f_{1}\left(\zeta ; x_{1}\right)\right)=0, & 0<x_{1}<\Lambda \\ \frac{\partial \delta u}{\partial x_{2}}\left(\zeta ; x_{1}, b\right)=T\left(\delta u\left(\zeta ; x_{1}, b\right)\right) & 0<x_{1}<\Lambda .\end{cases}
$$

Multiplying the differential equation in the adjoint problem (3.2) by $\overline{\delta u(\zeta ; \cdot)}$ and the differential equation in (3.16) by $u_{n}^{*}\left(\zeta_{;} \cdot\right)$, and integrating over the domain $D_{j}$ for $j=1, \cdots, \ell$, it follows that

$$
\begin{aligned}
& \int_{D_{j}}\left(\Delta u_{n}^{*}(\zeta ; \cdot)+k_{0}^{2} \overline{\varepsilon_{r}} u_{n}^{*}(\zeta ; \cdot)\right) \overline{\delta u(\zeta ; \cdot)}-u_{n}^{*}(\zeta ; \cdot) \overline{\left(\Delta \delta u(\zeta ; \cdot)+k_{0}^{2} \varepsilon_{r} \delta u(\zeta ; \cdot)\right)} d x \\
= & \int_{D_{j}} u_{n}^{*}(\zeta ; \cdot) k_{0}^{2} \overline{\delta \varepsilon_{r} u(\zeta ; \cdot)} d x, \quad j=1, \cdots, \ell .
\end{aligned}
$$

Applying the Green's formula on the left-hand side and adding all the equations together yields

$$
\begin{aligned}
& \int_{\Gamma_{j}(\zeta)}\left(\partial_{\nu} u_{n}^{*}(\zeta ; \cdot)\right)_{-} \overline{(\delta u(\zeta ; \cdot))_{-}}-\left(u_{n}^{*}(\zeta ; \cdot)\right)_{-} \overline{\left(\partial_{\nu} \delta u(\zeta ; \cdot)\right)_{-}} d s \\
& +\int_{\Gamma_{j}(\zeta)}\left(u_{n}^{*}(\zeta ; \cdot)\right)_{+} \overline{\left(\partial_{\nu} \delta u(\zeta ; \cdot)\right)_{+}}-\left(\partial_{\nu} u_{n}^{*}(\zeta ; \cdot)\right)_{+} \overline{(\delta u(\zeta ; \cdot))_{+}} d s \\
& +\int_{0}^{\Lambda} e^{i \kappa_{n} x} \overline{\delta u\left(\zeta ; x_{1}, b\right)} d x_{1}=k_{0}^{2} \int_{D} \overline{\delta \varepsilon_{r} u(\bar{\zeta} ; \cdot) u_{n}^{*}(\zeta ; \cdot) d x,}
\end{aligned}
$$

where we have used the boundary conditions in (3.2) and (3.16). By the continuity conditions along the interface $\Gamma_{j}(\zeta)$, this can be further reduced to the following:

$$
\begin{equation*}
\int_{0}^{\Lambda} e^{i \kappa_{n} x} \overline{\delta u\left(\zeta ; x_{1}, b\right)} d x_{1}=k_{0}^{2} \int_{D} \bar{\delta} \varepsilon_{r} u(\bar{\zeta} ; \cdot) u_{n}^{*}(\zeta ; \cdot) d x \tag{3.17}
\end{equation*}
$$

Now,

$$
\begin{aligned}
R^{\delta}(\zeta ; \boldsymbol{\alpha}) & =\sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}}\left|r_{n}+\delta r_{n}\right|^{2} \\
& =\sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}}\left\{\left|r_{n}\right|^{2}+2 \operatorname{Re}\left[r_{n} \overline{\delta r_{n}}\right]+\left|\delta r_{n}\right|^{2}\right\} .
\end{aligned}
$$

From the definition of (2.17), it follows that $\delta r_{n}=O(\delta u)$ where $\delta u$ is the perturbation of $u$. Due to the perturbation analysis of the boundary value problem (2.13), we have $\delta u=O\left(\delta \varepsilon_{r}\right)$ and $\delta \varepsilon_{r}=\left(\varepsilon_{r, 1}-\varepsilon_{r, 2}\right) \cdot \delta f_{j}$. Then it follows $\left|\delta r_{n}\right|^{2}=O\left(\left(\delta f_{j}\right)^{2}\right)$. We have

$$
R^{\delta}(\zeta ; \boldsymbol{\alpha})=R(\zeta ; \boldsymbol{\alpha})+2 \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[r_{n} \overline{\delta r_{n}}\right]+O\left(\left(\delta f_{j}\right)^{2}\right) .
$$

The perturbation $\delta R:=R^{\delta}-R$ is given by

$$
\begin{equation*}
\delta R=2 \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[r_{n} \overline{\delta r_{n}}\right]+O\left(\left(\delta f_{j}\right)^{2}\right) . \tag{3.18}
\end{equation*}
$$

For each term $r_{n} \overline{\delta r_{n}}$, by virtue of (2.17), it follows that

$$
r_{n} \overline{\delta r_{n}}= \begin{cases}\hat{u}_{n}(\zeta ; b) \cdot \frac{1}{\Lambda} \int_{0}^{\Lambda} e^{i \kappa_{n} x_{1}} \overline{\delta u\left(\zeta ; x_{1}, b\right)} d x_{1}, & n \neq 0  \tag{3.19}\\ \left(\hat{u}_{n}(\zeta ; b)-e^{-2 i k_{k} b+i p b}\right) \cdot \frac{1}{\Lambda} \int_{0}^{\Lambda} e^{i i_{n} x_{1}} \overline{\delta u\left(\zeta ; x_{1}, b\right)} d x_{1}, & n=0 .\end{cases}
$$

Therefore, substituting (3.17) into (3.19) yields

$$
r_{n} \overline{\delta r_{n}}= \begin{cases}\hat{u}_{n}(\zeta ; b) \cdot \frac{k_{0}^{2}}{\Lambda} \int_{D} \bar{\delta} \varepsilon_{r} \bar{u} u_{n}^{*} d x, & n \neq 0,  \tag{3.20}\\ \left(\hat{u}_{n}(\zeta ; b)-e^{-2 i k_{l} b+i \rho b}\right) \cdot \frac{k_{0}^{2}}{\Lambda} \int_{D} \bar{\delta} \varepsilon_{r} \bar{u} u_{n}^{*} d x, & n=0 .\end{cases}
$$

As such

$$
\delta R=\frac{2 k_{0}^{2}}{\Lambda} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[\left(\hat{u}_{n}(\zeta ; b)-\alpha_{n} e^{-2 i k_{e} b+i \rho b}\right) \cdot \int_{D} \bar{\delta}_{r} \bar{u} u_{n}^{*} d x\right]+O\left(\delta f_{j}^{2}\right) .
$$

Using (3.15), we arrive at

$$
\delta R=\frac{2 k_{0}^{2}}{\Lambda} \sum_{n \in \mathcal{N}} \frac{\eta_{n}}{\eta_{0}} \operatorname{Re}\left[\left.\left(\hat{u}_{n}(\zeta ; b)-\alpha_{n} e^{-2 i k_{k} b+i \rho b}\right) \cdot \overline{\left(\varepsilon_{r, 1}-\varepsilon_{r, 2}\right)} \cdot \int_{0}^{\Lambda}\left[\bar{u} u_{n}^{*}\right]\right|_{\left(x_{1}, f_{j}\right)} \cdot \delta f_{j} d x_{1}\right]+O\left(\delta f^{2}\right) .
$$

The desired formula (3.1b) for $D_{\alpha_{j}} R(j=2, \cdots, \ell)$ then follows by the chain rule.

## 4 The stochastic gradient descent method for the optimal design problems

### 4.1 The computational algorithm

To minimize $Q(\boldsymbol{\alpha})=E[R(\zeta ; \boldsymbol{\alpha})]$, the full gradient descent method applies the iteration

$$
\begin{equation*}
\boldsymbol{\alpha}^{(n+1)}=\boldsymbol{\alpha}^{(n)}-h_{n} \cdot D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right) \tag{4.1}
\end{equation*}
$$

where $h_{n}$ is the step length and $D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)$ is the gradient of $Q(\boldsymbol{\alpha})$ with respect to $\boldsymbol{\alpha}$. If the Monte Carlo method is used to sample the probability space, then

$$
D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)=E\left[D_{\alpha} R\left(\zeta ; \boldsymbol{\alpha}^{(n)}\right)\right] \approx \frac{1}{M_{C}} \sum_{m=1}^{M_{C}} D_{\alpha} R\left(\zeta_{m} ; \alpha^{(n)}\right)
$$

Usually the sampling size $M_{C}$ needs to be very large to obtain reasonably accurate approximation, and the computation of each $D_{\alpha} R\left(\zeta_{m} ; \alpha^{(n)}\right)$ requires solving the boundary value problem (2.13) and the adjoint problems (3.2). Therefore, computing the full gradient $D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)$ is very expensive during the iteration process.

Here we employ the stochastic gradient descent method to solve the optimization problems (I)-(III). The stochastic gradient descent method plays a significant role in solving large-scale modern machine learning problems and it is computationally efficient when the data set is large [7]. Its application for minimizing the objective function $Q(\boldsymbol{\alpha})$ is given by

$$
\begin{equation*}
\boldsymbol{\alpha}^{(n+1)}=\boldsymbol{\alpha}^{(n)}-h_{n} \cdot D_{\alpha} R\left(\zeta_{n} ; \boldsymbol{\alpha}^{(n)}\right) \tag{4.2}
\end{equation*}
$$

For each iteration $n$, the sample $\zeta_{n}$ is randomly chosen. In addition, at each iteration, the numerical method avoids the sampling of the gradient over the probability space and it requires the computation of the gradient $D_{\alpha} R\left(\zeta_{n}\right)$ only for one sample, although the convergence rate is slower than the full gradient algorithm above. The iterative sequence is not determined uniquely by the function $Q(\boldsymbol{\alpha})$, the starting point $\boldsymbol{\alpha}^{(1)}$, and the sequence of step size $\left\{h_{n}\right\}_{n=1}^{\infty}$. Rather, $\left\{\boldsymbol{\alpha}^{(n)}\right\}_{n=1}^{\infty}$ is a stochastic process whose behavior is determined by the random sequence $\left\{\zeta_{n}\right\}_{n=1}^{\infty}$.

The stochastic gradient descent method and the full gradient descent method offer different trade-offs in terms of computational cost at each iteration and the convergence rate for the iteration process. The full gradient iteration (4.1) is costly but stable, while the stochastic gradient descent iteration (4.2) is efficient but less stable. The mini-batch stochastic gradient descent method is designed to combine the advantages of both methods by choosing a small random samples of the gradients at each iteration. More precisely, the iteration takes the form

$$
\begin{equation*}
\boldsymbol{\alpha}^{(n+1)}=\boldsymbol{\alpha}^{(n)}-h_{n} \cdot G\left(\boldsymbol{\alpha}^{(n)}\right), \tag{4.3}
\end{equation*}
$$

where

$$
G\left(\boldsymbol{\alpha}^{(n)}\right)=\frac{1}{M_{0}} \sum_{m=1}^{M_{0}} D_{\alpha} R\left(\zeta_{m} ; \boldsymbol{\alpha}^{(n)}\right)
$$

is the average of the gradient over a small randomly chosen sample subset $\left\{\zeta_{m}\right\}_{m=1}^{M_{0}}$. The original stochastic gradient descent iteration (4.2) is a special case when $M_{0}=1$. When $M_{0}>1$, the mini-batch stochastic gradient descent method reduces the variance of the randomly chosen gradient during the iteration process by sampling over a larger set, thus it is more stable than the original algorithm (4.2). Here we apply the iteration formula (4.3) to solve the optimization problems (I)- (III) described as follows.

```
Algorithm 1 The mini-batch stochastic gradient descent method for Problems (I)- (III)
    Choose initial guess \(\boldsymbol{\alpha}^{(0)}\) and the sample size \(M_{0}\).
    while The average gradient of the sample set \(\left\|G\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}>\) tolerance do
```

        - Choose a random samples subset \(\left\{\zeta_{m}\right\}_{m=1}^{M_{0}}\).
        - For each sample \(\zeta_{m}\), solve the boundary value problem (2.13) and the adjoint
        problems (3.2) in Section 3.2.
            - Compute the gradient \(D_{\alpha} R\left(\zeta_{m} ; \cdot\right)\) by the formulas (3.1a) and (3.1b) in Section 3.2.
            - \(\operatorname{Set} G\left(\boldsymbol{\alpha}^{(n)}\right)=\frac{1}{M_{0}} \sum_{m=1}^{M_{0}} D_{\alpha} R\left(\zeta_{m} ; \boldsymbol{\alpha}^{(n)}\right)\).
            - Set \(\boldsymbol{\alpha}^{(n+1)}=\boldsymbol{\alpha}^{(n)}-h_{n} \cdot G\left(\boldsymbol{\alpha}^{(n)}\right), h_{n}>0\).
    end while
    
### 4.2 Convergence of the stochastic gradient descent algorithm

In this section, we examine the convergence of the stochastic gradient descent method. Let us focus on the case when the random variables $\left\{\tilde{\xi}_{j, m}\right\}_{m=0}^{M}$ in the Karhunen-Loève expansion (2.15) are uniformly distributed over the interval $[-0.5,0.5]$. In what follows, $C$ denotes a generic constant. Its value may vary from step to step but should be clear from the context.

Theorem 4.1. Assume the stochastic gradient descent iterations in Algorithm 1 satisfies the following conditions:
(1) The step sizes (learning rates) $\left\{h_{n}\right\}_{n=1}^{\infty}$ satisfy $\sum_{n=1}^{\infty} h_{n}=\infty$ and $\sum_{n=1}^{\infty} h_{n}^{2}<\infty$.
(2) The iteration sequence $\left\{\boldsymbol{\alpha}^{(n)}\right\}_{n=1}^{\infty}$ is bounded in the closed region $U_{\boldsymbol{\alpha}}=\left(\left[0, \beta_{1}\right] \times\left[\beta_{2}, \Lambda\right]\right)^{\ell}$ for certain constants $\beta_{1}, \beta_{2}>0$.

Then $E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2}\right] \rightarrow 0$ as $n \rightarrow \infty$.
To prove the theorem, we need the following lemma.
Lemma 4.2 ( [15], Lemma 1). Let $\left(a_{t}\right)_{t \geq 1},\left(b_{t}\right)_{t \geq 1}$ be two nonnegative real sequences. Assume that $\sum_{t=1}^{\infty} a_{t} b_{t}$ converges and $\sum_{t=1}^{\infty} a_{t}$ diverges, and there exists $K \geq 0$ such that $\mid b_{t+1}-$ $b_{t} \mid \leq K a_{t}$. Then $b_{t}$ converges to 0 .
Proof of Theorem 4.1 From the expression (2.15), we have

$$
\frac{\partial f_{j}\left(\zeta ; \alpha_{j} ; \cdot\right)}{\partial \alpha_{j}^{(1)}}=\bar{f}_{j}\left(\zeta ; \alpha_{j} ; \cdot\right),
$$

where $\bar{f}_{j}$ is given in (2.16). It follows that

$$
E\left[\left\|\frac{\partial f_{j}\left(\zeta ; \alpha_{j} ; x_{1}\right)}{\partial \alpha_{j}^{(1)}}\right\|_{L^{2}([0, \Lambda])}^{2}\right]=E\left[\left\|\bar{f}_{j}\left(\zeta ; \alpha_{j} ; x_{1}\right)\right\|_{L^{2}([0, \Lambda])}^{2}\right] \leq C \sum_{p=0}^{P}\left|\bar{\lambda}_{j p}\right|^{2} .
$$

On the other hand, $\frac{\partial f_{j}\left(\zeta ; \alpha_{j} ; x_{1}\right)}{\partial \alpha_{j}^{(1)}}=\alpha_{j}^{(1)} \cdot \frac{\partial \bar{f}_{j}\left(\zeta ; \alpha_{j}^{(2)} ; \cdot\right)}{\partial \alpha_{j}^{(2)}}$, and

$$
\begin{align*}
\frac{\partial \bar{f}_{j}\left(\zeta ; \alpha_{j}^{(2)} ; x_{1}\right)}{\partial \alpha_{j}^{(2)}}= & \frac{\bar{\lambda}_{j 0}^{\prime}}{2 \sqrt{\bar{\lambda}_{j 0}} \zeta_{0}(\zeta) \sqrt{\frac{1}{\Lambda}}+\sum_{p=1}^{P} \frac{\bar{\lambda}_{j p}^{\prime}}{2 \sqrt{\bar{\lambda}_{j p}}}\left[\zeta_{p}^{s}(\zeta) \sqrt{\frac{2}{\Lambda}} \sin \left(\frac{2 p \pi x_{1}}{\Lambda}\right)\right.} \\
& \left.+\zeta_{p}^{c}(\zeta) \sqrt{\frac{2}{\Lambda}} \cos \left(\frac{2 p \pi x_{1}}{\Lambda}\right)\right] \tag{4.4}
\end{align*}
$$

where

$$
\bar{\lambda}_{j p}^{\prime}:=\frac{\partial \bar{\lambda}_{j p}}{\partial \alpha_{j}^{(2)}}=\frac{\Lambda}{2} \frac{\partial \hat{c}_{j p}}{\partial \alpha_{j}^{(2)}} \quad \text { for } p=0,1,2, \cdots, P_{0}
$$

Recall that $\left\{\hat{c}_{j p}\right\}_{p=0}^{\infty}$ are the Fourier coefficients of the analytic function $\exp \left(-x_{1}^{2} /\left(\alpha_{j}^{(2)}\right)^{2}\right)$, thus $\left\{\frac{\partial \hat{c}_{j p}}{\partial \alpha_{j}^{(2)}}\right\}_{p=0}^{\infty}$ are the Fourier coefficients of the function $\frac{x_{1}^{2}}{\left(\alpha_{j}^{(2)}\right)^{3}} \exp \left(-x_{1}^{2} /\left(\alpha_{j}^{(2)}\right)^{2}\right)$,
which again is analytic for $x_{1} \in[0, \Lambda]$ and $\alpha_{j}^{(2)} \in\left[\beta_{2}, \Lambda\right]$. We deduce that

$$
E\left[\left\|\frac{\partial f_{j}\left(\zeta ; x_{1}\right)}{\partial \alpha_{j}^{(2)}}\right\|_{L^{2}([0, \Lambda])}^{2}\right] \leq C\left(\alpha_{j}^{(1)}\right)^{2} \sum_{p=0}^{P_{0}}\left|\frac{\lambda_{j p}^{\prime}}{2 \sqrt{\lambda_{j p}}}\right|^{2}=\frac{\Lambda^{2}}{8}\left(\alpha_{j}^{(1)}\right)^{2} \sum_{p=0}^{P_{0}}\left|\frac{\hat{c}_{j p}^{\prime}}{\sqrt{\hat{c}_{j p}}}\right|^{2} \leq C\left(P_{0}\right),
$$

where the constant $C$ depends on $P_{0}$. Namely, there holds

$$
\begin{equation*}
E\left[\left\|D_{\alpha_{j}} f_{j}\left(\zeta ; x_{1}\right)\right\|_{2}^{2}\right] \leq C, \quad j=1, \cdots, \ell . \tag{4.5}
\end{equation*}
$$

Let $p\left(y_{1}, \cdots, y_{2 P_{0}+1}\right)$ be the joint probability density function of the multivariate random variable $\left(\xi_{0}, \xi_{1}^{S}, \cdots, \xi_{P_{0}}^{S}, \xi_{1}^{c}, \cdots, \xi_{P_{0}}^{c}\right)$. From the continuous dependence of the solution to the boundary value problem (2.13) on the interfaces, $u$ can be viewed as a continuous function of random variables $\left(\xi_{0}, \xi_{1}^{s}, \cdots, \xi_{P_{0}}^{s}, \xi_{1}^{c}, \cdots, \xi_{P_{0}}^{c}\right)$. Therefore,

$$
\begin{align*}
& E\left[\left\|u\left(\xi_{0}, \xi_{1}^{s}, \cdots, \xi_{P_{0}}^{s}, \xi_{1}^{c}, \cdots, \xi_{P_{0}}^{c} ; \cdot\right)\right\|_{L^{2}(D)}^{2}\right] \\
= & \int_{\Omega}\left\|u\left(\xi_{0}, \xi_{1}^{s}, \cdots, \xi_{P_{0}}^{s}, \xi_{1}^{c}, \cdots, \xi_{P_{0}}^{c} ;\right)\right\|_{L^{2}(D)}^{2} d P\left(\xi_{0}, \xi_{1}^{s}, \cdots, \xi_{P_{0}}^{s}, \xi_{1}^{c}, \cdots, \xi_{P_{0}}^{c}\right)  \tag{4.6}\\
= & \int_{[-0.5,0.5]^{2 P_{0}+1}}\left\|u\left(y_{1}, \cdots, y_{2 P_{0}+1} ; \cdot\right)\right\|_{L^{2}(D)}^{2} p\left(y_{1}, \cdots, y_{2 P_{0}+1}\right) d y_{1} \cdots d y_{2 P_{0}+1} \leq C\left(P_{0}\right) .
\end{align*}
$$

Similarly, we have

$$
\begin{equation*}
E\left[\left\|u_{n}^{*}\right\|_{L^{2}(D)}^{2}\right] \leq C, E\left[\left\|\frac{\partial u}{\partial v}\right\|_{L^{2}(0, \Lambda)}^{2}\right] \leq C \text { and } E\left[\left\|\frac{\partial u_{n}^{*}}{\partial v}\right\|_{L^{2}(0, \Lambda)}^{2}\right] \leq C \tag{4.7}
\end{equation*}
$$

From Theorem 3.1,

$$
\begin{aligned}
E\left[\left\|D_{\alpha_{1}} R(\zeta ; \boldsymbol{\alpha})\right\|_{2}^{2}\right] & \leq C \max _{n \in \mathcal{N}} E\left[\left\|\hat{u}_{n}(\zeta ; b)+\int_{0}^{\Lambda}\left(\frac{\partial \bar{u}}{\partial v} \cdot \frac{\partial u_{n}^{*}}{\partial v} v_{2}\right) \cdot D_{\alpha_{1}} f_{1} d x_{1}\right\|_{2}^{2}\right] \\
& \leq C \max _{n \in \mathcal{N}}\left[E\left[\left|\hat{u}_{n}(\zeta ; b)\right|^{2}\right]+E\left[\left\|\int_{0}^{\Lambda}\left(\frac{\partial \bar{u}}{\partial v} \frac{\partial u_{n}^{*}}{\partial v} v_{2}\right) \cdot D_{\alpha_{1}} f_{1} d x_{1}\right\|_{2}^{2}\right]\right]
\end{aligned}
$$

Using (4.5)-(4.7) and the Cauchy-Schwartz inequality, we obtain

$$
E\left[\left\|D_{\alpha_{1}} R(\zeta ; \alpha)\right\|_{2}^{2}\right] \leq C
$$

Similarly, we have

$$
E\left[\left\|D_{\alpha_{j}} R(\zeta ; \alpha)\right\|_{2}^{2}\right] \leq C, j=2, \cdots, \ell
$$

Following the same lines, it can also be shown that

$$
\begin{equation*}
\left\|D_{\alpha}^{2} Q(\boldsymbol{\alpha})\right\|_{2}^{2}=E\left[\left\|D_{\alpha}^{2} R(\boldsymbol{\alpha})\right\|_{2}^{2}\right] \leq C \tag{4.8}
\end{equation*}
$$

Here $D_{\alpha}^{2}$ denotes the Hessian matrix of $Q(\boldsymbol{\alpha})$.

Without loss of generality, we assume that $M_{0}=1$ in Algorithm 1, and the iteration becomes $\boldsymbol{\alpha}^{(n+1)}=\boldsymbol{\alpha}^{(n)}-h_{n} \cdot D_{\alpha} R\left(\zeta ; \boldsymbol{\alpha}^{(n)}\right)$. Now, there exists $\widetilde{\boldsymbol{\alpha}} \in\left(\boldsymbol{\alpha}^{(n)}, \boldsymbol{\alpha}^{(n+1)}\right)$ such that

$$
\begin{align*}
& Q\left(\boldsymbol{\alpha}^{(n+1)}\right)-Q\left(\boldsymbol{\alpha}^{(n)}\right) \\
\leq & D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)^{T}\left(\boldsymbol{\alpha}^{(n+1)}-\boldsymbol{\alpha}^{(n)}\right)+\frac{1}{2}\left\|D_{\boldsymbol{\alpha}}^{2} Q(\widetilde{\boldsymbol{\alpha}})\right\|_{2}^{2} \cdot\left\|\boldsymbol{\alpha}^{(n+1)}-\boldsymbol{\alpha}^{(n)}\right\|_{2}^{2} \\
\leq & D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)^{T}\left(\boldsymbol{\alpha}^{(n+1)}-\boldsymbol{\alpha}^{(n)}\right)+\frac{1}{2} C\left\|\boldsymbol{\alpha}^{(n+1)}-\boldsymbol{\alpha}^{(n)}\right\|_{2}^{2}  \tag{4.9}\\
= & -h_{n} D_{\boldsymbol{\alpha}} Q\left(\boldsymbol{\alpha}^{(n)}\right)^{T} D_{\alpha} R\left(\zeta_{n} ; \boldsymbol{\alpha}^{(n)}\right)+\frac{1}{2} h_{n}^{2} C\left\|D_{\alpha} R\left(\zeta_{n} ; \boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2} .
\end{align*}
$$

Let us take conditional expectation of (4.9) with respect to $\zeta_{n}$. Then $E_{\zeta_{n}}\left[\left\|D_{\alpha} R\left(\zeta_{n} ; \boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2} \leq\right.$ $C$ since $E\left[\left\|D_{\alpha} R\left(\zeta_{n} ; \boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2} \leq C\right.$. It follows that

$$
\begin{aligned}
E_{\zeta_{n}}\left[Q\left(\boldsymbol{\alpha}^{(n+1)}\right)\right]-Q\left(\boldsymbol{\alpha}^{(n)}\right) & \leq-h_{n} D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)^{T} E_{\zeta_{n}}\left[D_{\alpha} R\left(\zeta_{n} ; \boldsymbol{\alpha}^{(n)}\right)\right]+\frac{1}{2} h_{n}^{2} C^{2} \\
& \leq-h_{n}\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2}+\frac{1}{2} h_{n}^{2} C^{2} .
\end{aligned}
$$

Taking the expectation yields

$$
\begin{equation*}
E\left[Q\left(\boldsymbol{\alpha}^{(n+1)}\right)\right]-E\left[Q\left(\boldsymbol{\alpha}^{(n)}\right)\right] \leq-h_{n} E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2}\right]+\frac{1}{2} h_{n}^{2} C^{2} . \tag{4.10}
\end{equation*}
$$

Denote $Q_{\infty}:=\liminf _{n \rightarrow \infty} Q\left(\boldsymbol{\alpha}^{(n)}\right)$, then $Q_{\infty}>0$. By adding (4.10) from 1 to $n$, we obtain

$$
Q_{\infty}-E\left[Q\left(\boldsymbol{\alpha}^{(1)}\right)\right] \leq E\left[Q\left(\boldsymbol{\alpha}^{(n+1)}\right)\right]-E\left[Q\left(\boldsymbol{\alpha}^{(1)}\right] \leq-\sum_{j=1}^{n} h_{j} E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(j)}\right)\right\|_{2}^{2}\right]+\frac{1}{2} C^{2} \sum_{j=1}^{n} h_{j}^{2} .\right.
$$

Hence,

$$
\sum_{j=1}^{n} h_{j} E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2}\right] \leq E\left[Q\left(\boldsymbol{\alpha}^{(1)}\right)\right]-Q_{\infty}+\frac{1}{2} C^{2} \sum_{j=1}^{n} h_{j}^{2} .
$$

Since $\sum_{n=1}^{\infty} h_{n}^{2}<\infty$, there holds

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \sum_{j=1}^{n} h_{j} E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(j)}\right)\right\|_{2}^{2}\right]<\infty . \tag{4.11}
\end{equation*}
$$

On the other hand, since $U_{\alpha}$ is closed, it follows that

$$
\begin{align*}
E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n+1)}\right)\right\|_{2}^{2}\right]-E\left[\left\|D_{\alpha} Q\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2}\right] & \leq \max _{\alpha \in U_{\alpha}}\left\|D_{\alpha}^{2} Q\right\|_{2}^{2} \cdot E\left[\boldsymbol{\alpha}^{(n+1)}-\boldsymbol{\alpha}^{(n)}\right] \\
& \leq h_{n} \max _{\alpha \in U_{\alpha}}\left\|D_{\alpha}^{2} Q\right\|_{2}^{2} \cdot E\left[D_{\alpha} R\left(\zeta_{n}, \boldsymbol{\alpha}^{(n)}\right)\right]  \tag{4.12}\\
& \leq C^{2} h_{n} .
\end{align*}
$$

From (4.11),(4.12) and Lemma 4.2, we deduce that $E\left[\left\|D_{\boldsymbol{\alpha}} Q\left(\boldsymbol{\alpha}^{(n)}\right)\right\|_{2}^{2}\right] \rightarrow 0$ as $n \rightarrow \infty$.

## 5 Numerical experiments

In this section, we present several numerical examples to demonstrate the efficiency of the numerical algorithm for solving the optimal design problems. The first numerical example tests the efficiency of the stochastic gradient descent method for solving the optimization problem and its performance compared to the full gradient descent method. In the second example, we apply Algorithm 1 to solve the optimization problem (I) at fixed frequency. Example 3 and Example 4 demonstrate the efficiency of Algorithm 1 when it is applied to the optimization problem (II) and (III) in the case of broad-band frequency and multiple incident angles. In all examples, we set the sample size $M_{0}=5$ and use $\tilde{R}:=\frac{1}{M_{0}} \sum_{m=1}^{M_{0}} R\left(\zeta_{m} ; \boldsymbol{\alpha}^{(n)}\right)$ and $\|\tilde{G}\|:=\frac{1}{M_{0}}\left\|\sum_{m=1}^{M_{0}} D_{\alpha} R\left(\zeta_{m} ; \boldsymbol{\alpha}^{(n)}\right)\right\|_{2}$ to denote the average reflectivity and the average amplitude of the gradient at each iteration. For all examples, the average thickness of the each layer is set as 300 nm , and the size of the periodic cell $\Lambda=1500 \mathrm{~nm}$. The state and adjoint problems are solved by the finite element method, where for each random realization of the optical structure, each subdomain $D_{j}$ is discretized by a triangular mesh such that the interface $\Gamma_{j}$ shares the common side of two adjacent triangles located above and below $\Gamma_{j}$ respectively.

Example 1 For simplicity we do not explicitly consider the glass substrate and assume that the solar cell consists of an absorbing layer (e.g., a-Si:H) at the bottom and a transparent conducting oxide (TCO) layer on the top. The bottom of the structure $\Gamma_{1}$, and the interface $\Gamma_{2}$ between the absorbing layer and the TCO layer are randomly texutured. Assume that the free space wavelength $\lambda_{0}=650 \mathrm{~nm}$. The relative permittivity of the TCO layer is $\varepsilon_{r, 1}=3.667$, and the relative permittivity for the absorbing layer is $\varepsilon_{r, 2}=17.6380+0.3780 i[10,13,18]$. We consider the configuration when the incident angle $\theta=0$.

Assume that the interfaces $\Gamma_{1}$ and $\Gamma_{2}$ are random processes with the covariance function $c_{j}\left(x_{1}-\tilde{x}_{1}\right)=\left(\alpha_{j}^{(1)}\right)^{2} \exp \left(-\left|x_{1}-\tilde{x}_{1}\right|^{2} /\left(\alpha_{j}^{(2)}\right)^{2}\right)$. The initial guess is chosen to be $\left(\alpha_{1}, \alpha_{2}\right)=(35 \mathrm{~nm}, 20 \mathrm{~nm})$. We apply both the full gradient descent method (4.1) and the stochastic gradient descent method described in Algorithm 1 to solve the optimization problem (2.19), where the gradient $D_{\alpha} R\left(\zeta_{m} ; \cdot\right)$ is computed via formulas (3.1a) and (3.1b). The stopping criteria is set as the amplitude of the average gradient amplitude $\|\tilde{G}\|$ being less than 0.05 .

Figure 3 shows the value of the average reflectivity $\tilde{R}$ at each iteration for the stochastic gradient method when the random variables in the Karhunen-Loève expansion (2.14) are chosen to be uniform and Gaussian random variables, respectively. For the former, the reflectivity $\tilde{R}$ decreases quickly in the first 50 iterations and it takes about 180 iterations to achieve the stopping criteria, while it only takes about 60 iterations for the latter to achieve the same tolerance. For completeness we also show the amplitude of the gradient $\|\tilde{G}\|$ at each iteration in Figure 4. It is clear that while $\|\tilde{G}\|$ oscillates during the
iterations but the envelope of $\|\tilde{G}\|$ decreases as $n$ increases. This is consistent with our convergence analysis presented in Section 4.

Figure 5 shows the reflectivity $Q$ for each iteration when the full gradient descent method is applied. Here the Monte Carlo method is used for sampling the random space. Table 1 and Table 2 collect the optimal parameters obtained by two different numerical approaches. It is observed that the optimal parameters obtained by the stochastic gradient descent method and the full gradient descent method are close to each other.


Figure 3: The reflectivity value $\tilde{R}(\boldsymbol{\alpha})$ during the stochastic gradient descent iterations. Top: $\left\{\tilde{\xi}_{0}, \xi_{m, s}, \xi_{m, c}\right\}$ are uniform random variables; Bottom: $\left\{\tilde{\xi}_{0}, \xi_{m, s}, \xi_{m, c}\right\}$ are Gaussian random variables.

Table 1: The optimal values of $\alpha^{(1)}$ and $\alpha^{(2)}$ obtained by the full gradient method and the stochastic gradient method for uniform random variables.

| optimal result(nm) | $\alpha_{1}^{(1)}$ | $\alpha_{1}^{(2)}$ | $\alpha_{2}^{(1)}$ | $\alpha_{2}^{(2)}$ | reflectivity |
| :---: | :---: | :---: | :---: | :---: | :---: |
| Full gradient method | 41 | 30 | 38 | 26 | 0.503 |
| Stochastic gradient method | 44 | 33 | 35 | 24 | 0.495 |

Though the optimal solutions obtained by the full gradient and the stochastic gradient methods are close, their computational cost is significantly different. When the Monte Carlo method is applied to sample the probability space, the sample size needs to be large. In our numerical experiment, the sample size is chosen to be 1000 for each iteration. When the optical wavelength is 650 nm , the computation of the gradient $D_{\alpha} R(\zeta, \alpha)$ for each sample $\zeta$ requires solving the boundary value problem (2.13) once and 9 adjoint problems (3.2) with all propagating modes. The full gradient algorithm stops after 7 steps, thus it requires solving 70000 boundary value problems. On the other hand, the 60 stochastic gradient descent iterations only requires solving no more than 3000 boundary


Figure 4: The amplitude of the average gradient $\|\tilde{G}\|$ during the stochastic gradient descent iterations. Top: $\left\{\xi_{0}, \xi_{m, s}, \xi_{m, c}\right\}$ are uniform random variables; Bottom: $\left\{\xi_{0}, \xi_{m, s}, \xi_{m, c}\right\}$ are Gaussian random variables.


Figure 5: The reflectivity value $Q(\boldsymbol{\alpha})$ during the full gradient iterations. Top: $\left\{\xi_{0}, \xi_{m, s}, \xi_{m, c}\right\}$ are uniform random variables; Bottom: $\left\{\tilde{\xi}_{0}, \xi_{m, s}, \xi_{m, c}\right\}$ are Gaussian random variables.

Table 2: The optimal values of $\alpha^{(1)}$ and $\alpha^{(2)}$ obtained by the full gradient method and the stochastic gradient method for Gaussian random variables

| optimal result(nm) | $\alpha_{1}^{(1)}$ | $\alpha_{1}^{(2)}$ | $\alpha_{2}^{(1)}$ | $\alpha_{2}^{(2)}$ | reflectivity |
| :---: | :---: | :---: | :---: | :---: | :---: |
| Full gradient method | 55 | 67 | 40 | 17 | 0.364 |
| Stochastic gradient method | 57 | 63 | 42 | 15 | 0.352 |

value problems. Therefore, the stochastic gradient descent method lowers the computational cost dramatically.

Example 2 Consider the multiple-layer solar cell structure as shown in Figure 6, where the interfaces of the two TCO layers are patterened randomly. The refractive index of the TCO, the absorbing layer and the glass substrate are $1.915,4.2+0.045 i$ and 1.4 , respectively. Let the incident angle $\theta=0$ and the wavelength $\lambda_{0}=650 \mathrm{~nm}$. We assume that all the interfaces are Gaussian random processes.

It takes about 120 iterations for the stochastic gradient descent method to achieve the desired tolerance, and the average reflectivity $\tilde{R}$ for each iteration is shown in Figure 7. The reflectivity $\tilde{R}$ decreases from the initial value 0.75 to 0.37 for the optimal random structure, with the corresponding absorptance value 0.63 . As a comparison, the absorptance of the structure with all flat interfaces is only 0.13 . Figure 8 depicts the wave field for one realization of random structure with the optimal result and Figure 9 shows the wave field for the structure with flat interfaces. It is observed that the waves are scattered in the random medium, which elongates the optical path and increases the overall absorptance of the structure.


Figure 6: Optical structure with four random interfaces.


Figure 7: The reflectivity value $\tilde{R}(\boldsymbol{\alpha})$ during the stochastic gradient iterations for Problem(I). The multi-layered medium has four random interfaces shown in Figure 6.


Figure 8: Numerical solution of the boundary value problem (2.13) for one realization of random structure with optimal parameters.



Figure 9: Numerical solution of the boundary value problem (2.13) when all the interfaces are flat.

Example 3 In this example, we consider the more challenging optimization problem with multiple frequencies, which is formulated in (2.20). Let us still use the multi-layer structure shown in Figure 6. Assume that the interfaces are Gaussian random processes and the incident angle $\theta=0$. The refractive index of the absorbing layer is set as $4.5+0.12 i$ and $4.2+0.045 i$ when $\lambda_{\text {min }}=600 \mathrm{~nm}$ and $\lambda_{\text {max }}=650 \mathrm{~nm}$, respectively. For simplicity, we assume the refractive index of the absorbing layer is a linear function of the wavelength between $\lambda_{\text {min }}$ and $\lambda_{\text {max }}$.

The integral $\int_{\lambda_{\text {min }}}^{\lambda_{\text {max }}} R(\zeta ; \boldsymbol{\alpha}, \lambda) d \lambda$ in (2.20) is approximated by the $\operatorname{sum} \frac{1}{M_{\lambda}} \sum_{m=1}^{M_{\lambda}} R\left(\zeta ; \boldsymbol{\alpha}, \lambda_{m}\right)$, in which $\lambda_{m}=\lambda_{\min }+\frac{\lambda_{\max }-\lambda_{\text {min }}}{M_{\lambda}-1}(m-1)$ for $m=1, \cdots, M_{\lambda}$. We consider the normal incidence with the incident angle $\theta=0$.

Figure 10 shows the average reflectivity at each iteration for the stochastic gradient descent approach. It is calculated the average reflectivity of the optimal structure is about 0.28 and the average absorptance is about 0.72 . We see that the absorptance is significantly enhanced compared to the structure with flat interfaces, which only attains the value 0.24 .
Example 4 In this example, we consider the optimization problem with multiple incident angle for the structure shown in Figure 6. The problem is formulated in (2.21). We still consider the interfaces with Gaussian random processes and solve the problem with the stochastic gradient descent method. In the calculation, the integral $\int_{\theta_{\min }}^{\theta_{\max }} R(\zeta ; \alpha, \theta) d \lambda$ is approximated with the sum $\frac{1}{M_{\theta}} \sum_{m=1}^{M_{\theta}} R\left(\zeta ; \alpha, \theta_{m}\right)$, where $\theta_{m}=\theta_{\min }+\frac{\theta_{\max }-\theta_{\min }}{M_{\theta}-1}(m-1)$ for


Figure 10: The reflectivity value $\tilde{R}(\boldsymbol{\alpha})$ during the stochastic gradient iterations for Problem (II)
$m=1, \cdots, M_{\theta}$. When the range of the incident angle starts from $\theta_{\min }=-\frac{\pi}{12}$ to $\theta_{\max }=\frac{\pi}{12}$, the reflectivity at each stochastic gradient iteration is shown in Figure 11. We obtain an average absorptance value of about 0.51 for the optimal structure, which again is significantly higher than the the structure with flat interface with an absorptance value 0.11 .


Figure 11: The reflectivity value $\tilde{R}(\boldsymbol{\alpha})$ during the stochastic gradient iterations for Problem (III).

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