# Analysis of an Inverse Problem Arising in Photolithography 

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#### Abstract

We consider the inverse problem of determining an optical mask that produces a desired circuit pattern in photolithography. We set the problem as a shape design problem in which the unknown is a two-dimensional domain. The relationship between the target shape and the unknown is modeled through diffractive optics. We develop a variational formulation that is well-posed and propose an approximation that can be shown to have convergence properties. The approximate problem can serve as a foundation to numerical methods.


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## 1 Introduction

Photolithography is a key process in the production of integrated circuits. It is the process by which circuit patterns are transferred onto silicon wafers. A review of this manufacturing technology is given in [16]. The main step in photolithography is the creation of a circuit image on the photoresist coating which sits on the silicon layer that is to be patterned. The image is formed using ultra-violet (UV) light which is diffracted by a mask, and refracted by a system of lenses. The mask simply consists of cut-outs, and lets light through the holes. The parts of the photoresist that are exposed to the UV light can be removed, leaving openings to the layer to be patterned. The next stage is etching, which removes material in the layer that is unprotected by the photoresist. Once etching is done, the photoresist can be removed, and the etched away "channels" may be filled. The entire process is illustrated schematically in Figure 1.

The problem we address in this work is the inverse problem of determining what mask is needed in order to remove a desired shape in the photoresist. The difficulty of producing a desired shape comes from the fact that the UV light is diffracted at the mask. Moreover, the chemicals in the photoresist reacts nonlinearly to UV exposure - only portions of the photoresist that have been exposed to a certain level of intensity are removed in the bleaching process.

The nature of the present work is analytical. Our goal is to formulate mathematically well-posed problems for photolithography. The methods we use to prove well-posedness are constructive and may serve as a foundation for a computational method.

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Figure 1: The photolithographic process. Ultraviolet light, diffracted by a mask, forms an image on the photoresist. The exposed portion of the photoresist is removed, leaving openings. Etching removes parts of the layer to be patterned. After etching, the photoresist is removed.

Our investigation into photolithography is inspired by the work of Cobb [3] who was the first to approach this problem from the point of view of optimal design which utilizes a physically-based model. This general approach was further developed by introducing a level set method in [18]. A different computational approach which models the mask as a pixelated binary image can be found in [14].

The plan of the paper is as follows. In the first and preliminary section, Section 2, we develop the most basic model for removal of the exposed photoresist. We describe the inverse problem to be solved. This is followed by a discussion of the approximate problem whose properties we intend to investigate in this work. Section 3 contains mathematical preliminaries needed for our work. We introduce the basic notation and recall various results which will be useful for our analysis. In particular, in Subsection 3.3 we discuss the geometry of masks or circuits and how to measure the distance between two of them. In Section 4 we discuss the properties of the operator which maps the mask into the circuit. Section 5 provides an analysis of the variational approach to the problem of the optimization of the mask and we prove a convergence result for it, Theorem 5.5, in the framework of $\Gamma$ convergence.

## 2 Description of the inverse problem

This section is separated into three subsections. First, we review some basic facts about Fourier transforms and prove a result about approximation of a Gaussian. We follow this with a discussion of the optics involved and a model for photolithography. In the final subsection we describe the inverse problem and its approximation.

### 2.1 Fourier transform and approximation of Gaussians

We first set some notation and describe a few preliminary results. For every $x \in \mathbb{R}^{2}$, we shall set $x=\left(x_{1}, x_{2}\right)$, where $x_{1}$ and $x_{2} \in \mathbb{R}$. For every $x \in \mathbb{R}^{2}$ and $r>0$, we shall denote by $B_{r}(x)$ the open ball in $\mathbb{R}^{2}$ centered at $x$ of radius $r$. Usually we shall write $B_{r}$ instead of $B_{r}(0)$. We recall that, for any set $E \subset \mathbb{R}^{2}$, we denote by $\chi_{E}$ its characteristic function, and for any $r>0, B_{r}(E)=\bigcup_{x \in E} B_{r}(x)$.

For any $f \in \mathcal{S}^{\prime}\left(\mathbb{R}^{2}\right)$, the space of tempered distributions, we denote by $\hat{f}$ its Fourier transform, which, if $f \in L^{1}\left(\mathbb{R}^{2}\right)$, may be written as

$$
\hat{f}(\xi)=\int_{\mathbb{R}^{2}} f(x) \mathrm{e}^{-\mathrm{i} \xi \cdot x} \mathrm{~d} x, \quad \xi \in \mathbb{R}^{2} .
$$

We recall that $f(x)=(2 \pi)^{-2} \hat{\hat{f}}(-x)$, that is, when also $\hat{f} \in L^{1}\left(\mathbb{R}^{2}\right)$,

$$
f(x)=\frac{1}{(2 \pi)^{2}} \int_{\mathbb{R}^{2}} \hat{f}(\xi) \mathrm{e}^{\mathrm{i} \xi \cdot x} \mathrm{~d} \xi, \quad x \in \mathbb{R}^{2} .
$$

If $f$ is a radial function, that is $f(x)=\phi(|x|)$ for any $x \in \mathbb{R}^{2}$, then

$$
\hat{f}(\xi)=2 \pi \mathcal{H}_{0}(\phi)(|\xi|), \quad \xi \in \mathbb{R}^{2}
$$

where

$$
\mathcal{H}_{0}(\phi)(s)=\int_{0}^{+\infty} r J_{0}(s r) \phi(r) \mathrm{d} r, \quad s \geq 0
$$

is the Hankel transform of order $0, J_{0}$ being the Bessel function of order 0 , see for instance [4].

We denote the Gaussian distribution by $G(x)=(2 \pi)^{-1} \mathrm{e}^{-|x|^{2} / 2}, x \in \mathbb{R}^{2}$, and let us note that $\hat{G}(\xi)=\mathrm{e}^{-|\xi|^{2} / 2}, \xi \in \mathbb{R}^{2}$. Moreover, $\|G\|_{L^{1}\left(\mathbb{R}^{2}\right)}=1$. Furthermore if $\delta_{0}$ denotes the Dirac delta centered at 0 , we have $\widehat{\delta_{0}} \equiv 1$, therefore $(2 \pi)^{-2} \hat{1}=\delta_{0}$.

For any function $f$ defined on $\mathbb{R}^{2}$ and any positive constant $s$, we denote $f_{s}(x)=$ $s^{-2} f(x / s), x \in \mathbb{R}^{2}$. We note that $\left\|f_{s}\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}=\|f\|_{L^{1}\left(\mathbb{R}^{2}\right)}$ and $\widehat{f}_{s}(\xi)=\hat{f}(s \xi), \xi \in \mathbb{R}^{2}$.

We conclude these preliminaries with the following integrability result for the Fourier transform and its applications.

Theorem 2.1 There exists an absolute constant $C$ such that the following estimate holds

$$
\|\hat{f}\|_{L^{1}\left(\mathbb{R}^{2}\right)} \leq C\|f\|_{W^{2,1}\left(\mathbb{R}^{2}\right)} .
$$

Proof. This result is contained in Theorem A in [9] and it is based on previous analysis done in [13].

We recall that a more detailed analysis on conditions for which integrability of the Fourier transform holds may be found in [17]. However the previous result is simple to use and it is enough for our purposes, in particular for proving the following lemma.

Lemma 2.2 For any $\tilde{\delta}>0$ there exist a constant $s_{0}, 0<s_{0} \leq 1$, and a radial function $\hat{T} \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$ such that $\hat{T} \equiv 1$ on $B_{s_{0}}$ and, if we call $T=(2 \pi)^{-2} \hat{T}$, then $T \in W^{2,1}\left(\mathbb{R}^{2}\right)$ and

$$
\|T-G\|_{W^{1,1}\left(\mathbb{R}^{2}\right)} \leq \tilde{\delta}
$$

Proof. We sketch the proof of this result. Let us consider the following cut-off function $\phi \in C^{\infty}(\mathbb{R})$ such that $\phi$ is nonincreasing, $\phi \equiv 1$ on $(-\infty, 0]$ and $\phi \equiv 0$ on $[1,+\infty)$.

We define a function $\hat{\tilde{T}}$ as follows

$$
\hat{\tilde{T}}(x)=\phi(|x|-1)+(1-\phi(|x|-1)) \widehat{G_{s_{0}}}(x) \phi(|x|-b), \quad x \in \mathbb{R}^{2},
$$

for suitable constants $s_{0}, 0<s_{0} \leq 1$ and $b \geq 2$. We call $\tilde{T}=(2 \pi)^{-2} \hat{\tilde{T}}$.
Then lengthy but straightforward computations, with the aid of Theorem 2.1, allow us to prove that for some $s_{0}$ small enough and for some $b=s_{0}^{-1} b_{0}$, with $b_{0}$ large enough, we have

$$
\left\|\tilde{T}-G_{s_{0}}\right\|_{W^{1,1}\left(\mathbb{R}^{2}\right)} \leq \tilde{\delta}
$$

Then, let $\hat{T}(x)=\hat{\tilde{T}}\left(x / s_{0}\right), x \in \mathbb{R}^{2}$, so that $T=\tilde{T}_{1 / s_{0}}$, or equivalently $\tilde{T}=T_{s_{0}}$. Therefore

$$
\left\|T_{s_{0}}-G_{s_{0}}\right\|_{W^{1,1}\left(\mathbb{R}^{2}\right)} \leq \tilde{\delta}
$$

By a simple rescaling argument we have that $\hat{T}$ satisfies the required properties. Furthermore, by this construction, we may choose $\hat{T}$ such that it is radially nonincreasing, $\hat{T} \equiv 1$ on $B_{s_{0}}$ and it decays to zero in a suitable smooth, exponential way.

### 2.2 A model of image formation

We are now in the position to describe the model we shall use. The current industry standard for modeling the optics is based on Kirchhoff approximation. Under this approximation, the light source at the mask is on where the mask is open, and off otherwise (see Figure 1). Propagation through the lenses can be calculated using Fourier optics. It is further assumed that the image plane, in this case the plane of the photoresist, is at the focal distance of the optical system. If there were no diffraction, a perfect image of the mask would be formed on the image plane. Diffraction, together with partial coherence of the light source, acts to distort the formed image.

The mask, which we mention consists of cut-outs, is represented as a binary function, i.e., it is a characteristic function of the cut-outs. Suppose that $D$ represents the cut-outs, then the mask is given by

$$
m(x)=\chi_{D}(x) .
$$

The image is the light intensity on the image plane. This is given by [12]

$$
\begin{equation*}
I(x)=\int_{\mathbb{R}^{2}} \int_{\mathbb{R}^{2}} m(\xi) K(x-\xi) J(\xi-\eta) K(x-\eta) m(\eta) \mathrm{d} \xi \mathrm{~d} \eta, \quad x \in \mathbb{R}^{2} . \tag{2.1}
\end{equation*}
$$

In the above expression the kernel $K(\cdot)$ is called the coherent point spread function and describes the optical system. For an optical system with a circular aperture, once the wavenumber of the light used, $k>0$, has been chosen, the kernel depends on a single parameter called the Numerical Aperture, NA. Notice that the wavelength is $\lambda=2 \pi / k$. Let us recall that the so-called Jinc function is defined as

$$
\operatorname{Jinc}(x)=\frac{J_{1}(|x|)}{2 \pi|x|}, \quad x \in \mathbb{R}^{2}
$$

where $J_{1}$ is the Bessel function of order 1. We notice that in the Fourier space, see for instance [6, page 14],

$$
\widehat{\operatorname{Jinc}}(\xi)=\chi_{B_{1}}(\xi), \quad \xi \in \mathbb{R}^{2}
$$

If we denote by $s=(k \mathrm{NA})^{-1}$, then the kernel is usually modeled as follows

$$
K(x)=\operatorname{Jinc}_{s}(x)=\frac{k \mathrm{NA}}{2 \pi} \frac{J_{1}(k \mathrm{NA}|x|)}{|x|}, \quad x \in \mathbb{R}^{2},
$$

therefore

$$
\hat{K}(\xi)=\chi_{B_{1}}(s \xi)=\chi_{B_{1 / s}}(s \xi)=\chi_{B_{k \mathrm{NA}}}(\xi), \quad \xi \in \mathbb{R}^{2}
$$

If NA goes to $+\infty$, that is $s \rightarrow 0^{+}$, then $\hat{K}$ converges pointwise to 1 , thus $K$ approximates in a suitable sense the Dirac delta.

For technical reasons, we shall consider a slightly different coherent point spread function $K$. Let us fix a positive constant $\tilde{\delta}$, to be chosen later. We shall replace the characteristic function $\chi_{B_{1}}$, the Fourier transform of the Jinc function, with the function $\hat{T}\left(s_{0} \xi\right), \xi \in \mathbb{R}^{2}$, with $\hat{T}$ and $s_{0}$ as in Lemma 2.2. Therefore $\hat{T}\left(s_{0} \cdot\right)$ is a radial function that it is still identically equal to 1 on $B_{1}$, it is still compactly supported, it is nonincreasing with respect to the radial variable and it decays to zero in a smooth, exponential way. Its Fourier transform is $T_{s_{0}}$ and we shall assume that

$$
\begin{equation*}
K(x)=\left(T_{s_{0}}\right)_{s}(x)=T_{s s_{0}}(x), \quad x \in \mathbb{R}^{2} \tag{2.2}
\end{equation*}
$$

where again $s=(k \mathrm{NA})^{-1}$. Also in this model, if NA goes to $+\infty$, that is $s \rightarrow 0^{+}$, then $\hat{K}$ converges pointwise to 1 , thus $K$ approximates in a suitable sense the Dirac delta.

The function $J(\cdot)$ is called the mutual intensity function. If the illumination is fully coherent, $J \equiv 1$. In practice, illumination is never fully coherent and is parametrized by a coherency coefficient $\sigma$. A typical model for $J$ is

$$
\begin{equation*}
J(x)=2 \frac{J_{1}(k \sigma \mathrm{NA}|x|)}{k \sigma \mathrm{NA}|x|}=\pi \operatorname{Jinc}(k \sigma \mathrm{NA}|x|), \quad x \in \mathbb{R}^{2} . \tag{2.3}
\end{equation*}
$$

Thus,

$$
\frac{1}{(2 \pi)^{2}} \hat{J}(\xi)=\frac{1}{\pi(k \sigma \mathrm{NA})^{2}} \chi_{B_{k \sigma \mathrm{NA}}}(\xi), \quad \xi \in \mathbb{R}^{2},
$$

that, as $\sigma \rightarrow 0^{+}$, converges, in a suitable sense, to the Dirac delta. Therefore full coherence is achieved for $\sigma \rightarrow 0^{+}$. In fact, if $\sigma \rightarrow 0^{+}, J$ converges to 1 uniformly on any compact subset of $\mathbb{R}^{2}$. The equation (2.1) is often referred to as the Hopkins areal intensity representation. As it will become apparent from the analysis developed in the paper, the value of $s$ is related to the scale of details that the manufacturing of the mask allows, thus in turn to the scale of details of the desired circuit. Therefore, we typically consider $k \mathrm{NA} \gg 1$, that is $s \ll 1$, and $k \sigma \mathrm{NA} \ll 1$.

### 2.3 The inverse problem and its approximation

The photoresist material responds to the intensity of the image. When intensity at the photoresist goes over a certain threshold, it is then considered exposed and can be removed. Therefore, the exposed pattern, given a mask $m(x)$, is

$$
\begin{equation*}
\Omega=\left\{x \in \mathbb{R}^{2}: I(x)>h\right\}, \tag{2.4}
\end{equation*}
$$

where $h$ is the exposure threshold. Clearly, $\Omega$ depends on the mask function $m(x)$, which we recall is given by the characteristic function of $D$ representing the cut-outs, that is $\Omega=\Omega(D)$. In photolithography, we have a desired exposed pattern which we wish to achieve. The inverse problem is to find a mask that achieves this desired exposed pattern. Mathematically, this cannot, in general, be done. Therefore, the inverse problem must be posed as an optimal design problem.

Suppose the desired pattern is given by $\Omega_{0}$. We pose the minimization problem

$$
\begin{equation*}
\min _{D \in \mathcal{A}} d\left(\Omega(D), \Omega_{0}\right) . \tag{2.5}
\end{equation*}
$$

The distance function $d(\cdot, \cdot)$ will be discussed in detail below. The admissible set $\mathcal{A}$ is our search space, and needs to be defined carefully as well.

Instead of solving (2.5), we pose a variational problem for a function $u$ (instead of the mask $D$ ). We will show below that this problem is well-posed and that as the approximation parameter is set to zero, we recover the solution of (2.5) under a perimeter penalization.

Instead of dealing with the characteristic function $\chi_{D}(x)$ which represents the mask, we will work with a phase-field function $u$ which takes on values of 0 and 1 with smooth transitions. Thus, the intensity in (2.4) is calculated with $u$ instead of $m=\chi_{D}$ in (2.1), so $I$ is a function of $u$. At this point, we will not be precise about the space of functions to which $u$ belongs. To force $u$ to take on values of mostly 0 and 1 , we introduce the Mordica-Mortola energy

$$
P_{\varepsilon}(u)=\frac{1}{\varepsilon} \int W(u)+\varepsilon \int|\nabla u|^{2},
$$

where $W(t)=9 t^{2}(t-1)^{2}$ is a double-well potential. We will regularize the problem of minimizing the distance between the target pattern and the exposed region by this energy.

Then we relax the hard threshold in defining the exposed region $\Omega$ in (2.4). Let $\phi(t)$ be a $C^{\infty}$ nondecreasing approximate Heaviside function with values $\phi(t \leq-1 / 2)=0$ and $\phi(t \geq 1 / 2)=1$. The function

$$
\Phi_{\eta}(u)=\phi\left(\frac{I(u)-h}{\eta}\right)
$$

will be 1 where the intensity $I \geq h+\eta / 2$. A sigmoidal threshold function is employed in the computational work in [14].

Now we consider the distance function between $\Omega$ and $\Omega_{0}$ in (2.5). Let

$$
\begin{equation*}
d=d\left(\Omega, \Omega_{0}\right)=\int\left|\chi_{\Omega}-\chi_{\Omega_{0}}\right|+\left|P(\Omega)-P\left(\Omega_{0}\right)\right|, \tag{2.6}
\end{equation*}
$$

where $\chi_{\Omega}$ is the characteristic function of the set $\Omega$ and $P(\Omega)$ is the perimeter of the region $\Omega$. To approximate this distance function, we replace it by

$$
d_{\eta}\left(u, \Omega_{0}\right)=\int\left|\Phi_{\eta}(u)-\chi_{\Omega_{0}}\right|+\left|\int\right| \nabla\left(\Phi_{\eta}(u)\right)\left|-P\left(\Omega_{0}\right)\right| .
$$

The characteristic function of $\Omega$ is replaced by the smooth threshold function while its perimeter is replaced by the TV-norm of the function.

The approximate problem we shall solve is

$$
F_{\varepsilon}(u)=d_{\eta(\varepsilon)}\left(u, \Omega_{0}\right)+b P_{\varepsilon}(u) \rightarrow \min .
$$

The remainder of the paper is an analytical study of this minimization problem. We will show that it is well-posed, and that in the limit $\varepsilon \rightarrow 0^{+}$, we recover the solution of the original problem (2.5) under a perimeter penalization.

## 3 Mathematical preliminaries

By $\mathcal{H}^{1}$ we denote the 1-dimensional Hausdorff measure and by $\mathcal{L}^{2}$ we denote the 2-dimensional Lebesgue measure. We recall that, if $\gamma \subset \mathbb{R}^{2}$ is a smooth curve, then $\mathcal{H}^{1}$ restricted to $\gamma$ coincides with its arclength. For any Borel $E \subset \mathbb{R}^{2}$ we denote $|E|=\mathcal{L}^{2}(E)$.

Let $\mathcal{D}$ be a bounded open set contained in $\mathbb{R}^{2}$, with boundary $\partial \mathcal{D}$. We say that $\mathcal{D}$ has a Lipschitz boundary if for every $x=\left(x_{1}, x_{2}\right) \in \partial \mathcal{D}$ there exist a Lipschitz function $\varphi: \mathbb{R} \rightarrow \mathbb{R}$ and a positive constant $r$ such that for any $y \in B_{r}(x)$ we have, up to a rigid transformation,

$$
y=\left(y_{1}, y_{2}\right) \in \mathcal{D} \quad \text { if and only if } \quad y_{2}<\varphi\left(y_{1}\right) .
$$

We note that $\mathcal{D}$ has a finite number of connected components, whereas $\partial \mathcal{D}$ is formed by a finite number of rectifiable Jordan curves, therefore $\mathcal{H}^{1}(\partial \mathcal{D})=$ length $(\partial \mathcal{D})<+\infty$.

We recall some basic notation and properties of functions of bounded variation and sets of finite perimeter. For a more comprehensive treatment of these subjects see, for instance, [1, 7, 8.

Given a bounded open set $\mathcal{D} \subset \mathbb{R}^{2}$, we denote by $B V(\mathcal{D})$ the Banach space of functions of bounded variation. We recall that $u \in B V(\mathcal{D})$ if and only if $u \in L^{1}(\mathcal{D})$ and its distributional derivative $D u$ is a bounded vector measure. We endow $B V(\mathcal{D})$ with the standard norm as follows. Given $u \in B V(\mathcal{D})$, we denote by $|D u|$ the total variation of its distributional derivative and we set $\|u\|_{B V(\mathcal{D})}=\|u\|_{L^{1}(\mathcal{D})}+|D u|(\mathcal{D})$. We shall call $P(u, \mathcal{D})=|D u|(\mathcal{D})$. We recall that whenever $u \in W^{1,1}(\mathcal{D})$, then $u \in B V(\mathcal{D})$ and $|D u|(\mathcal{D})=\int_{\mathcal{D}}|\nabla u|$, therefore $\|u\|_{B V(\mathcal{D})}=\|u\|_{L^{1}(\mathcal{D})}+\|\nabla u\|_{L^{1}(\mathcal{D})}=\|u\|_{W^{1,1}(\mathcal{D})}$.

We say that a sequence of $B V(\mathcal{D})$ functions $\left\{u_{h}\right\}_{h=1}^{\infty}$ weakly* converges in $B V(\mathcal{D})$ to $u \in B V(\mathcal{D})$ if and only if $u_{h}$ converges to $u$ in $L^{1}(\mathcal{D})$ and $D u_{h}$ weakly* converges to $D u$ in $\mathcal{D}$, that is

$$
\begin{equation*}
\lim _{h} \int_{\mathcal{D}} v \mathrm{~d} D u_{h}=\int_{\mathcal{D}} v \mathrm{~d} D u \quad \text { for any } v \in C_{0}(\mathcal{D}) . \tag{3.1}
\end{equation*}
$$

By Proposition 3.13 in [1], we have that if a sequence of $B V(\mathcal{D})$ functions $\left\{u_{h}\right\}_{h=1}^{\infty}$ is bounded in $B V(\mathcal{D})$ and converges to $u$ in $L^{1}(\mathcal{D})$, then $u \in B V(\mathcal{D})$ and $u_{h}$ converges to $u$ weakly* in $B V(\mathcal{D})$.

We say that a sequence of $B V(\mathcal{D})$ functions $\left\{u_{h}\right\}_{h=1}^{\infty}$ strictly converges in $B V(\mathcal{D})$ to $u \in B V(\mathcal{D})$ if and only if $u_{h}$ converges to $u$ in $L^{1}(\mathcal{D})$ and $\left|D u_{h}\right|(\mathcal{D})$ converges to $|D u|(\mathcal{D})$. Indeed,

$$
d_{s t}(u, v)=\int_{\mathcal{D}}|u-v|+||D u|(\mathcal{D})-|D v|(\mathcal{D})|
$$

is a distance on $B V(\mathcal{D})$ inducing the strict convergence. We also note that strict convergence implies weak ${ }^{*}$ convergence.

Let $\mathcal{D}$ be a bounded open set with Lipschitz boundary. A sequence of $B V(\mathcal{D})$ functions $\left\{u_{h}\right\}_{h=1}^{\infty}$ such that $\sup _{h}\left\|u_{h}\right\|_{B V(\mathcal{D})}<+\infty$ admits a subsequence converging weakly* in $B V(\mathcal{D})$ to a function $u \in B V(\mathcal{D})$, see for instance Theorem 3.23 in [1]. As a corollary, we
infer that for any $C>0$ the set $\left\{u \in B V(\mathcal{D}):\|u\|_{B V(\mathcal{D})} \leq C\right\}$ is a compact subset of $L^{1}(\mathcal{D})$.

For any fixed constant $R>0$, with a slight abuse of notation, we shall identify $L^{1}\left(B_{R}\right)$ with the set $\left\{u \in L^{1}\left(\mathbb{R}^{2}\right): u=0\right.$ a.e. outside $\left.B_{R}\right\}$.

Let $E$ be a bounded Borel set contained in $B_{R} \subset \mathbb{R}^{2}$. We shall denote by $\chi_{E}$ its characteristic function. We notice that $E$ is compactly contained in $B_{R+1}$, which we shall denote by $E \Subset B_{R+1}$. We say that $E$ is a set of finite perimeter if $\chi_{E}$ belongs to $B V\left(B_{R+1}\right)$ and we call the number $P(E)=\left|D \chi_{E}\right|\left(B_{R+1}\right)$ its perimeter. Analogously, for any $u \in$ $L^{1}\left(B_{R}\right) \cap B V\left(B_{R+1}\right)$, we shall denote $P\left(u, B_{R+1}\right)=|D u|\left(B_{R+1}\right)$. Obviously, if $u=\chi_{E}$, then $P\left(u, B_{R+1}\right)=P(E)$.

Let us further remark that the intersection of two sets of finite perimeter is still a set of finite perimeter. Moreover, whenever $E$ is open and $\mathcal{H}^{1}(\partial E)$ is finite, then $E$ is a set of finite perimeter, see for instance [7, Section 5.11, Theorem 1]. Therefore a bounded open set $\mathcal{D}$ with Lipschitz boundary is a set of finite perimeter and its perimeter $P(\mathcal{D})$ coincides with $\mathcal{H}^{1}(\partial \mathcal{D})$.

## 3.1 $\Gamma$-convergence approximation of the perimeter functional

Let us introduce the following, slightly different, version of a $\Gamma$-convergence result due to Modica and Mortola, [1]. We shall follow the notation and proofs contained in [2]. We begin by setting some notation. For the definition and properties of $\Gamma$-convergence we refer to [5].

For any bounded open set $\mathcal{D} \subset \mathbb{R}^{2}$, with a slight abuse of notation, we identify $W_{0}^{1, p}(\mathcal{D})$, $1<p<+\infty$, with the subset of $W^{1, p}\left(\mathbb{R}^{2}\right)$ functions $u$ such that $u$ restricted to $\mathcal{D}$ belongs to $W_{0}^{1, p}(\mathcal{D})$ and $u$ is equal to 0 almost everywhere outside $\mathcal{D}$. Let us assume that for some positive constant $R$ we have $\mathcal{D} \subset B_{R}$. We recall that any function in $L^{1}(\mathcal{D})$ is extended to zero outside $\mathcal{D}$ and the same procedure is used for $L^{1}\left(B_{R}\right)$. Therefore, with this slight abuse of notation, $L^{1}(\mathcal{D}) \subset L^{1}\left(B_{R}\right)$. Throughout the paper, for any $p, 1 \leq p \leq+\infty$, we shall denote its conjugate exponent by $p^{\prime}$, that is $p^{-1}+\left(p^{\prime}\right)^{-1}=1$.

Theorem 3.1 Let $\mathcal{D} \subset B_{R} \subset \mathbb{R}^{2}$ be a bounded open set with Lipschitz boundary. Let us also assume that $\mathcal{D}$ is convex.

Let $1<p<+\infty$ and $W: \mathbb{R} \rightarrow[0,+\infty)$ be a continuous function such that $W(t)=0$ if and only if $t \in\{0,1\}$. Let $c_{p}=\left(\int_{0}^{1}(W(s))^{1 / p^{\prime}} \mathrm{d} s\right)^{-1}$.

For any $\varepsilon>0$ we define the functional $P_{\varepsilon}: L^{1}\left(\mathbb{R}^{2}\right) \rightarrow[0,+\infty]$ as follows

$$
P_{\varepsilon}(u)= \begin{cases}\frac{c_{p}}{p^{\prime} \varepsilon} \int_{\mathcal{D}} W(u)+\frac{c_{p} \varepsilon^{p-1}}{p} \int_{\mathcal{D}}|\nabla u|^{p} & \text { if } u \in W_{0}^{1, p}(\mathcal{D})  \tag{3.2}\\ +\infty & \text { otherwise }\end{cases}
$$

Let $P: L^{1}\left(\mathbb{R}^{2}\right) \rightarrow[0,+\infty]$ be such that

$$
P(u)= \begin{cases}P\left(u, B_{R+1}\right) & \text { if } u \in B V\left(B_{R+1}\right), u \in\{0,1\} \text { a.e. }  \tag{3.3}\\ & \text { and } u=0 \text { a.e. outside } \mathcal{D} \\ +\infty & \text { otherwise }\end{cases}
$$

Then $P=\Gamma-\lim _{\varepsilon \rightarrow 0^{+}} P_{\varepsilon}$ with respect to the $L^{1}\left(\mathbb{R}^{2}\right)$ norm.

Remark 3.2 We observe that $P(u)=P(E)$ if $u=\chi_{E}$ where $E$ is a set of finite perimeter contained in $\overline{\mathcal{D}}$ and $P(u)=+\infty$ otherwise.

Furthermore, we note that the result does not change if in the definition of $P_{\varepsilon}$ we set $P_{\varepsilon}(u)=+\infty$ whenever $u$ does not satisfy the constraint

$$
\begin{equation*}
0 \leq u \leq 1 \text { a.e. in } \mathcal{D} . \tag{3.4}
\end{equation*}
$$

Proof. We sketch the proof following that of Theorem 4.13 in [2]. In fact, the only difference with respect to that theorem is that we assume $\mathcal{D}$ convex and that we take $W_{0}^{1, p}(\mathcal{D})$ instead of $W^{1, p}(\mathcal{D})$ in the definition of $P_{\varepsilon}$.

By Proposition 4.3 in [2], we obtain that $P(u) \leq \Gamma-\lim \inf _{\varepsilon \rightarrow 0^{+}} P_{\varepsilon}(u)$ for any $u \in L^{1}\left(\mathbb{R}^{2}\right)$. In order to obtain the $\Gamma$-lim sup inequality, we follow the procedure described in Section 4.2 of [2]. It would be enough to construct $\mathcal{M} \subset L^{1}\left(\mathbb{R}^{2}\right)$ such that the following two conditions are satisfied. First, we require that, for any $u \in L^{1}\left(\mathbb{R}^{2}\right)$ such that $P(u)<+\infty$, there exists a sequence $\left\{u_{j}\right\}_{j=1}^{\infty}$ such that $u_{j} \in \mathcal{M}$, for any $j \in \mathbb{N}, u_{j} \rightarrow u$ in $L^{1}\left(\mathbb{R}^{2}\right)$ as $j \rightarrow \infty$, and $P(u)=\lim _{j} P\left(u_{j}\right)$. Second, for any $u \in \mathcal{M}, \Gamma-\lim \sup _{\varepsilon \rightarrow 0^{+}} P_{\varepsilon}(u) \leq P(u)$.

We choose $\mathcal{M}=\left\{u=\chi_{E}: E \Subset \mathcal{D}, E\right.$ of class $\left.C^{\infty}\right\}$. The second property follows by Proposition 4.10 in [2]. As far as the first property is concerned, this can be obtained by following the proof of Theorem 1.24 in [8]. That theorem states that any bounded set of finite perimeter $E$ can be approximated by a sequence of $C^{\infty}$ sets $\left\{E_{j}\right\}_{j=1}^{\infty}$ such that, as $j \rightarrow \infty, \int_{\mathbb{R}^{2}}\left|\chi_{E_{j}}-\chi_{E}\right| \rightarrow 0$ and $P\left(E_{j}\right) \rightarrow P(E)$. If we assume that $E \subset \overline{\mathcal{D}}$, and that $\mathcal{D}$ is convex, by choosing in the proof of Theorem 1.24 in [8] a value of $t$ satisfying $1 / 2<t<1$, we obtain that the sets $E_{j}$ are also compactly contained in $\mathcal{D}$, for any $j \in \mathbb{N}$.

Also the following result, due to Modica, [10], will be useful.
Proposition 3.3 For any $C>0$, let us take $1<p<+\infty$ and any $\varepsilon>0$, and let us define

$$
A_{C}=\left\{u \in L^{1}\left(\mathbb{R}^{2}\right): 0 \leq u \leq 1 \text { a.e. and } P_{\varepsilon}(u) \leq C\right\} .
$$

Then $A_{C}$ is precompact in $L^{1}\left(\mathbb{R}^{2}\right)$.
Proof. We repeat, for the reader's convenience, the arguments developed in 10. Clearly $A_{C}$ is a bounded subset of $L^{1}(\mathcal{D})$. Let $\left\{u_{n}\right\}_{n=1}^{\infty}$ be a sequence in $A_{C}$. We need to prove that there exists a subsequence converging in $L^{1}(\mathcal{D})$. For any $t, 0 \leq t \leq 1$, let $\phi(t)=$ $\int_{0}^{t}(W(s))^{1 / p^{\prime}} \mathrm{d} s$. For any $n \in \mathbb{N}$, we define $v_{n}=\phi\left(u_{n}\right)$ and we observe that $0 \leq v_{n} \leq \phi(1)$ almost everywhere. Therefore, the functions $v_{n}, n \in \mathbb{N}$, are uniformly bounded in $L^{\infty}(\mathcal{D})$ and, consequently, in $L^{1}(\mathcal{D})$. Furthermore, since $\phi$ is a $C^{1}$ function, with bounded $C^{1}$ norm, then $D v_{n}=\phi^{\prime}\left(u_{n}\right) D u_{n}=W^{1 / p^{\prime}}\left(u_{n}\right) D u_{n}$. Therefore,

$$
\int_{\mathcal{D}}\left|D v_{n}\right|=\int_{\mathcal{D}}\left|W^{1 / p^{\prime}}\left(u_{n}\right)\right|\left|D u_{n}\right| \leq P_{\varepsilon}\left(u_{n}\right) / c_{p}
$$

We infer that there exists a subsequence $\left\{v_{n_{k}}\right\}_{k=1}^{\infty}$ converging, as $k \rightarrow \infty$, to a function $v_{0}$ in $L^{1}(\mathcal{D})$ and almost everywhere. Let $\psi$ be the inverse function of $\phi$ and let $u_{0}=\psi\left(v_{0}\right)$. We observe that $\psi$ is bounded and uniformly continuous on $[0, \phi(1)]$, hence we conclude that, as $k \rightarrow \infty, u_{n_{k}}$ converges to $u_{0}$ in $L^{1}(\mathcal{D})$.

Remark 3.4 With the same proof, we can show the following. Let us consider any family $\left\{u_{\varepsilon}\right\}_{0<\varepsilon \leq \varepsilon_{0}}$ such that, for some positive constant $C$ and for any $\varepsilon, 0<\varepsilon \leq \varepsilon_{0}$, we have $0 \leq u_{\varepsilon} \leq 1$ almost everywhere and $P_{\varepsilon}\left(u_{\varepsilon}\right) \leq C$. Then $\left\{u_{\varepsilon}\right\}_{0<\varepsilon \leq \varepsilon_{0}}$ is precompact in $L^{1}\left(\mathbb{R}^{2}\right)$.

### 3.2 Convolutions

We recall that, for any two functions $f$ and $g$ defined on $\mathbb{R}^{2}$, we define the convolution of $f$ and $g, f * g$, as follows

$$
(f * g)(x)=\int_{\mathbb{R}^{2}} f(x-y) g(y) \mathrm{d} y=\int_{\mathbb{R}^{2}} f(y) g(x-y) \mathrm{d} y, \quad x \in \mathbb{R}^{2}
$$

whenever this is well-defined.
The following classical properties of convolutions will be used. First convolution is commutative. Second, as a consequence of Young inequality we have the following result about integrability and regularity of convolutions.

Proposition 3.5 Let $1 \leq r, p, q \leq+\infty$ be such that $1+\frac{1}{r}=\frac{1}{p}+\frac{1}{q}$, and let $n=0,1,2, \ldots$.
Let $1 \leq q<+\infty$, let $f \in L^{q}\left(\mathbb{R}^{2}\right)$ and let $g \in W^{n, p}\left(\mathbb{R}^{2}\right)$. Then $h=f * g \in W^{n, r}\left(\mathbb{R}^{2}\right)$ and there exists a constant $C$, depending on $n, p, q$ and $r$ only, such that

$$
\|h\|_{W^{n, r}\left(\mathbb{R}^{2}\right)} \leq C\|f\|_{L^{q}\left(\mathbb{R}^{2}\right)}\|g\|_{W^{n, p}\left(\mathbb{R}^{2}\right)}
$$

Let $q=+\infty$ and let $f \in L^{\infty}\left(\mathbb{R}^{2}\right)$, with compact support. If $g \in W^{n, 1}\left(\mathbb{R}^{2}\right)$, then $h=$ $f * g \in W^{n, \infty}\left(\mathbb{R}^{2}\right)$ and there exists a constant $C$, depending on $n$ only, such that

$$
\|h\|_{W^{n, \infty}\left(\mathbb{R}^{2}\right)} \leq C\|f\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}\|g\|_{W^{n, 1}\left(\mathbb{R}^{2}\right)}
$$

If $f \in L^{1}\left(\mathbb{R}^{2}\right)$ and $g \in L^{\infty}\left(\mathbb{R}^{2}\right)$, then $h=f * g \in L^{\infty}\left(\mathbb{R}^{2}\right)$ and it holds $\|h\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq$ $\|f\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}\|g\|_{L^{1}\left(\mathbb{R}^{2}\right)}$. Furthermore, if $g$ is uniformly continuous and $\omega_{g}$ denotes its modulus of continuity, then $h$ is also uniformly continuous and

$$
\omega_{h} \leq\|f\|_{L^{1}\left(\mathbb{R}^{2}\right)} \omega_{g}
$$

Finally, let $f \in L^{1}\left(\mathbb{R}^{2}\right)$ and let $g \in C^{n, \alpha}\left(\mathbb{R}^{2}\right)$, for some $\alpha, 0<\alpha \leq 1$. Then $h \in C^{n, \alpha}\left(\mathbb{R}^{2}\right)$ and there exists a constant $C$, depending on $n$ and $\alpha$ only, such that

$$
\|h\|_{C^{n, \alpha}\left(\mathbb{R}^{2}\right)} \leq C\|f\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|g\|_{C^{n, \alpha}\left(\mathbb{R}^{2}\right)} .
$$

### 3.3 The geometry of masks and circuits

In this subsection we investigate the following two questions, namely what are reasonable assumptions on the geometry of the mask $D$ and how to measure the distance between the constructed circuit $\Omega$ and the desired one $\Omega_{0}$. We begin with the following definition. During this subsection, in most cases proofs will be omitted and left to the reader.

For given positive constants $r$ and $L$, we say that a bounded open set $\Omega \subset \mathbb{R}^{2}$ is Lipschitz or $C^{0,1}$ with constants $r$ and $L$ if for every $x \in \partial \Omega$ there exists a Lipschitz function $\varphi: \mathbb{R} \rightarrow \mathbb{R}$, with Lipschitz constant bounded by $L$, such that for any $y \in B_{r}(x)$, and up to a rigid transformation,

$$
\begin{equation*}
y=\left(y_{1}, y_{2}\right) \in \Omega \quad \text { if and only if } \quad y_{2}<\varphi\left(y_{1}\right) . \tag{3.5}
\end{equation*}
$$

Without loss of generality, we may always assume that $x=(0,0)$ and $\varphi(0)=0$. We shall always denote by $e_{1}$ and $e_{2}$ the vectors of the canonical bases. Clearly the orientation of the canonical bases may vary depending on $x \in \partial \Omega$.

We shall also use the following notation. There exist positive constants $\delta_{1} \leq 1 / 2, \delta_{2} \leq \delta_{1}$ and $m_{1} \leq 1$, all of them depending on $L$ only, such that the following holds. For any $x \in \partial \Omega$ and for any $\delta>0$, let $M_{\delta}(x)=\left\{y:\left|y_{1}\right| \leq \delta r, y_{2}=\varphi\left(y_{1}\right)\right\}$ and $N_{\delta}(x)=\left\{y:\left|y_{1}\right| \leq\right.$ $\left.\delta_{1} r, \varphi\left(y_{1}\right)-\delta r \leq y_{2} \leq \varphi\left(y_{1}\right)+\delta r\right\}$. Then we assume that, for any $\delta, 0<\delta \leq \delta_{2}$, the following properties hold. First, $N_{\delta}(x) \subset B_{r / 2}(x)$ (hence $M_{\delta_{1}}(x) \subset B_{r / 2}(x)$ as well). Clearly $N_{\delta}(x)$ is contained in $\bar{B}_{\delta r}(\partial \Omega)$, and we assume that $N_{\delta}(x)$ contains $\bar{B}_{m_{1} \delta r}\left(M_{\delta_{1} / 2}(x)\right)$ and that for any $y \in\left\{y:\left|y_{1}\right| \leq \delta_{1} r / 2, y_{2}=\varphi\left(y_{1}\right) \pm \delta r\right\}, y \notin \bar{B}_{m_{1} \delta r}(\partial \Omega)$.

For any integer $k=1,2, \ldots$, any $\alpha, 0<\alpha \leq 1$, and any positive constants $r$ and $L$, we say that a bounded open set $\Omega \subset \mathbb{R}^{2}$ is $C^{k, \alpha}$ with constants $r$ and $L$ if for every $x \in \partial \Omega$ there exists a $C^{k, \alpha}$ function $\varphi: \mathbb{R} \rightarrow \mathbb{R}$, with $C^{k, \alpha}$ norm bounded by $L$, such that for any $y \in B_{r}(x)$, and up to a rigid transformation, (3.5) holds. Without loss of generality, we may always assume that $x=(0,0)$ and $\varphi(0)=0$.

Let us fix three positive constants $r, L$ and $R$. Let $\mathcal{A}^{0,1}(r, L, R)$ be the class of all bounded open sets, contained in $B_{R} \subset \mathbb{R}^{2}$, which are Lipschitz with constants $r$ and $L$. For any integer $k=1,2, \ldots$ and any $\alpha, 0<\alpha \leq 1$, we denote with $\mathcal{A}^{k, \alpha}(r, L, R)$ the class of all bounded open sets, contained in $B_{R} \subset \mathbb{R}^{2}$, which are $C^{k, \alpha}$ with constants $r$ and $L$.

Since we shall identify open sets $D$ with their characteristic functions $\chi_{D}$, if $\mathcal{A}=$ $\mathcal{A}^{0,1}(r, L, R)$, (or $\mathcal{A}=\mathcal{A}^{k, \alpha}(r, L, R)$, respectively) then, with a slight abuse of notation, $\mathcal{A}$ will also denote the subset of functions $u \in L^{1}\left(B_{R}\right)$ such that $u=\chi_{D}$ for some $D \in \mathcal{A}$. Moreover, we shall denote

$$
A=\left\{u \in L^{1}\left(B_{R}\right): 0 \leq u \leq 1 \text { a.e. in } B_{R}\right\}
$$

and, for any $\gamma>0$,

$$
\begin{equation*}
\mathcal{A}_{\gamma}=\left\{u \in A:\left\|u-\chi_{D}\right\|_{L^{1}\left(B_{R}\right)} \leq \gamma \text { for some } D \in \mathcal{A}\right\} . \tag{3.6}
\end{equation*}
$$

Let us assume that $\Omega_{1}$ and $\Omega_{2}$ belong to $\mathcal{A}^{0,1}(r, L, R)$. There are several ways to define the distance between these two sets. We shall describe four of them and study their relationships. We let

$$
\begin{align*}
d_{1} & =d_{1}\left(\Omega_{1}, \Omega_{2}\right)=d_{H}\left(\bar{\Omega}_{1}, \bar{\Omega}_{2}\right) ;  \tag{3.7}\\
\tilde{d}_{1} & =\tilde{d}_{1}\left(\Omega_{1}, \Omega_{2}\right)=d_{H}\left(\partial \Omega_{1}, \partial \Omega_{2}\right) ;  \tag{3.8}\\
d_{2} & =d_{2}\left(\Omega_{1}, \Omega_{2}\right)=\left|\Omega_{1} \Delta \Omega_{2}\right|=\left\|\chi_{\Omega_{1}}-\chi_{\Omega_{2}}\right\|_{L^{1}\left(B_{R+1}\right)} ;  \tag{3.9}\\
d_{3} & =d_{3}\left(\Omega_{1}, \Omega_{2}\right)=d_{2}+\left|P\left(\Omega_{1}\right)-P\left(\Omega_{2}\right)\right|=d_{s t}\left(\chi_{\Omega_{1}}, \chi_{\Omega_{2}}\right) . \tag{3.10}
\end{align*}
$$

Here $d_{H}$ denotes the Hausdorff distance, whereas we recall that $P(\Omega)$ denotes the perimeter of $\Omega$ in $B_{R+1}$ and $d_{s t}$ is the distance inducing strict convergence in $B V\left(B_{R+1}\right)$. First of all, we observe that all of these are distances. We now investigate their relationships.

We begin with the first two, $d_{1}$ and $\tilde{d}_{1}$, and we notice that

$$
\begin{equation*}
\text { if } d_{1} \leq r / 4, \text { then } d_{1} \leq \tilde{d}_{1} . \tag{3.11}
\end{equation*}
$$

There exists a constant $c, 0<c \leq 1$, depending on $L$ only, such that

$$
d_{1} \geq c \min \left\{r, \tilde{d}_{1}\right\}
$$

Therefore,

$$
\text { if } \tilde{d}_{1} \leq r \text {, then } \tilde{d}_{1} \leq C d_{1} \text {, }
$$

where $C=1 / c$. Furthermore, if $d_{1} \leq(c / 2) r$, then $\tilde{d}_{1}$ must be less than or equal to $r$, so

$$
\text { if } d_{1} \leq(c / 2) r \text {, then } \tilde{d}_{1} \leq C d_{1}
$$

Moreover, we can find a constant $c_{1}, 0<c_{1} \leq 1$, depending on $L$ only, such that

$$
\text { if } \tilde{d}_{1} \leq c_{1} r \text {, then } d_{1} \leq \tilde{d}_{1} \leq C d_{1}
$$

We conclude that we can find a constant $c_{1}, 0<c_{1} \leq 1$, and a constant $C \geq 1$, both depending on $L$ only, such that

$$
\begin{equation*}
\text { if either } d_{1} \leq c_{1} r \text { or } \tilde{d}_{1} \leq c_{1} r \text {, then } d_{1} \leq \tilde{d}_{1} \leq C d_{1} \tag{3.12}
\end{equation*}
$$

Since $d_{1}$ and $\tilde{d}_{1}$ are bounded by $2 R$, we also have

$$
\begin{equation*}
\text { if both } d_{1} \geq c_{1} r \text { and } \tilde{d}_{1} \geq c_{1} r, \text { then } d_{1} \leq \frac{2 R}{c_{1} r} \tilde{d}_{1} \text { and } \tilde{d}_{1} \leq \frac{2 R}{c_{1} r} d_{1} . \tag{3.13}
\end{equation*}
$$

We finally observe that the estimates (3.12) and (3.13) are essentially optimal.
Before comparing $d_{1}$ (or $\tilde{d}_{1}$ ) with $d_{2}$ and $d_{3}$, let us make the following remark on the lengths of $\partial \Omega_{1}$ and $\partial \Omega_{2}$. If $\Omega$ is an open set which is Lipschitz with constants $r$ and $L$, then for any integer $n \geq 0$, we have

$$
\begin{equation*}
\mathcal{H}^{1}\left(\partial \Omega \cap\left(\bar{B}_{(n+1) r} \backslash B_{n r}\right)\right) \leq C(L) r(n+1) . \tag{3.14}
\end{equation*}
$$

Here, a simple computation shows that we may choose $C(L)=48 \sqrt{1+L^{2}}$.
Therefore, if we assume that $\Omega \subset B_{R}$ and $R \geq 10 r$, we may conclude that

$$
\begin{equation*}
P(\Omega) \leq C_{1}(L) R^{2} / r \tag{3.15}
\end{equation*}
$$

where $C_{1}(L)=\frac{1}{2}\left(\frac{11}{9}\right)^{2} C(L)$.
Moreover, there exist two constants $c_{2}, 0<c_{2} \leq c_{1}$, and $C_{1}>0$, depending on $L$ only, such that we have

$$
\begin{equation*}
\left|\bar{B}_{d}(\partial \Omega)\right| \leq C_{1} \text { length }(\partial \Omega) d \text { for any } d \leq c_{2} r . \tag{3.16}
\end{equation*}
$$

Since

$$
\text { if } \tilde{d}_{1} \leq c_{2} r \text {, then } d_{2} \leq \min \left\{\left|\bar{B}_{\tilde{d}_{1}}\left(\partial \Omega_{2}\right)\right|,\left|\bar{B}_{\tilde{d}_{1}}\left(\partial \Omega_{2}\right)\right|\right\},
$$

we obtain that

$$
\begin{equation*}
\text { if } \tilde{d}_{1} \leq c_{2} r \text {, then } d_{2} \leq C_{1} \min \left\{\operatorname{length}\left(\partial \Omega_{1}\right) \text {, length }\left(\partial \Omega_{2}\right)\right\} \tilde{d}_{1} \text {. } \tag{3.17}
\end{equation*}
$$

If $\tilde{d}_{1} \geq c_{2} r$, then $d_{2} \leq \pi R^{2} \leq \frac{\pi R^{2}}{c_{2} r} \tilde{d}_{1}$. By 3.15, we may conclude that

$$
\begin{equation*}
d_{2} \leq \frac{C_{2} R^{2}}{r} \tilde{d}_{1} \tag{3.18}
\end{equation*}
$$

Here $C_{2}$ depends on $L$ only. Moreover, up to changing the constants $c_{2}, C_{1}$ and $C_{2}$, (3.17) and (3.18) still hold if we replace $\tilde{d}_{1}$ with $d_{1}$.

On the other hand, there exists a constant $c_{3}, 0<c_{3} \leq \pi$, depending on $L$ only, such that

$$
d_{2} \geq c_{3} \min \left\{r^{2}, d_{1}^{2}\right\}
$$

We infer that either if $d_{1} \leq r$ or if $d_{2} \leq\left(c_{3} / 2\right) r^{2}$, then $d_{1} \leq C_{3} d_{2}^{1 / 2}$, where $C_{3}=1 / c_{3}^{1 / 2}$. If $d_{2} \geq\left(c_{3} / 2\right) r^{2}$, then $d_{1} \leq 2 R \leq \frac{4 C_{3}^{2}}{r^{2}} R d_{2}$ or, better, $d_{1} \leq \frac{2 \sqrt{2} C_{3}}{r} R d_{2}^{1 / 2}$. Summarizing, we have

$$
\begin{equation*}
\text { if } d_{2} \leq\left(c_{3} / 2\right) r^{2}, \text { then } d_{1} \leq C_{3} d_{2}^{1 / 2} \tag{3.19}
\end{equation*}
$$

and, finally,

$$
\begin{equation*}
\text { if } d_{2} \geq\left(c_{3} / 2\right) r^{2}, \text { then } d_{1} \leq \frac{2 \sqrt{2} C_{3}}{r} R d_{2}^{1 / 2} . \tag{3.20}
\end{equation*}
$$

Clearly, up to suitably changing the constants $c_{3}$ and $C_{3}$, the last two estimates still hold if we replace $d_{1}$ with $\tilde{d}_{1}$. We also remark that, as before, the estimates relating $d_{1}, \tilde{d}_{1}$ and $d_{2}$ are essentially optimal.

We have obtained that $d_{1}, \tilde{d}_{1}$ and $d_{2}$ are topologically equivalent distances. About $d_{2}$ and $d_{3}$, obviously $d_{2} \leq d_{3}$, however the two distances are not topologically equivalent. In fact we can find $\Omega$ and $\Omega_{i}, i \in \mathbb{N}$, open sets belonging to $\mathcal{A}^{0,1}(r, L, R)$, such that $d_{2}\left(\Omega, \Omega_{i}\right)$ goes to zero as $i \rightarrow \infty$, whereas $d_{3}\left(\Omega, \Omega_{i}\right) \geq c>0$ for any $i \in \mathbb{N}$. Therefore $d_{3}$ induces a strictly finer topology than the one induced by $d_{2}$

An assumption that the mask is a bounded open set which is Lipschitz with given constants $r$ and $L$ is reasonable from the manufacturing point of view as well as from the mathematical point of view, by the following compactness result.

Proposition 3.6 The set $\mathcal{A}^{0,1}(r, L, R)$ (respectively $\left.\mathcal{A}^{k, \alpha}(r, L, R), k=1,2, \ldots, 0<\alpha \leq 1\right)$ is compact with respect to the distance $d_{1}$.

We remark that the same result holds with respect to the distances $\tilde{d}_{1}$ and $d_{2}$. Furthermore, we obtain as a corollary that the set $\mathcal{A}_{\gamma}$ is closed with respect to the $L^{1}$ norm, for any $\gamma>0$.

The previous example shows that compactness fails with respect to the distance $d_{3}$, at least for the Lipschitz case. On the other hand, if $\Omega_{1}$ and $\Omega_{2}$ belong to $\mathcal{A}^{1, \alpha}(r, L, R)$, with $0<\alpha<1$, then, following Lemma 2.1 in [15], we can show that

$$
\begin{equation*}
\left|P\left(\Omega_{1}\right)-P\left(\Omega_{2}\right)\right| \leq C_{4}\left(\tilde{d}_{1}\left(\Omega_{1}, \Omega_{2}\right)\right)^{\alpha /(2 \alpha+2)}, \tag{3.21}
\end{equation*}
$$

where $C_{4}$ depends on $r, L, R$ and $\alpha$ only. We may conclude that in the $C^{k, \alpha}$ case, $k=1,2, \ldots$, $0<\alpha \leq 1, d_{3}$ is topologically equivalent to the other three distances and that Proposition 3.6 holds also with respect to the distance $d_{3}$.

It is worthwhile to observe that, under some circumstances, the estimate (3.21) can be extended to the piecewise $C^{1, \alpha}$ case. For example, typically we may assume that the desired circuit $\Omega_{0}$ belongs to $\mathcal{A}^{0,1}(r, L, R)$. Moreover, we assume that the boundary of $\Omega_{0}$ is composed by a finite number of closed segments $I_{i}, i=1, \ldots, n$, which are pairwise internally disjoint and whose lengths are greater than or equal to $2 r$. Therefore, $\Omega_{0}$ is actually a piecewise $C^{1, \alpha}$ open set. We shall show in Section 4 that, under suitable assumptions on the mask $D$, the corresponding constructed circuit $\Omega$ belongs to $\mathcal{A}^{1, \alpha}\left(r_{1}, L_{1}, \tilde{R}\right)$, for some suitable positive constants $r_{1} \leq r, L_{1} \geq L, \tilde{R} \geq R$ and $\alpha, 0<\alpha<1$. Then we can find
positive constants $c_{4}, 0<c_{4} \leq 1, C_{5}$ and $C_{6}$, depending on $r_{1}, L_{1}, \tilde{R}$ and $\alpha$ only, such that if $\tilde{d}_{1}\left(\Omega_{0}, \Omega\right) \leq c_{4} r_{1}$, then we can subdivide $\partial \Omega$ into smooth curves $J_{i}, i=1 \ldots, n$, which are pairwise internally disjoint, such that for any $i=1, \ldots, n$ we have

$$
d_{H}\left(J_{i}, I_{i}\right) \leq C_{5} \tilde{d}_{1}\left(\Omega_{0}, \Omega\right)
$$

and

$$
\text { length }\left(I_{i}\right)-2 C_{5} \tilde{d}_{1}\left(\Omega_{0}, \Omega\right) \leq \operatorname{length}\left(J_{i}\right) \leq \operatorname{length}\left(I_{i}\right)+C_{6}\left(\tilde{d}_{1}\left(\Omega_{0}, \Omega\right)\right)^{\alpha /(2 \alpha+2)}
$$

Therefore,

$$
-2 n C_{5} \tilde{d}_{1}\left(\Omega_{0}, \Omega\right) \leq P(\Omega)-P\left(\Omega_{0}\right) \leq n C_{6}\left(\tilde{d}_{1}\left(\Omega_{0}, \Omega\right)\right)^{\alpha /(2 \alpha+2)} .
$$

By these reasonings it might seem that we may choose to measure the distance between the desired circuit $\Omega_{0}$ and the reconstructed one $\Omega$ by using any of these distances. However, there are several reasons to prefer the distance $d_{3}$, which we actually choose. In fact, it is easier to compute than $d_{1}$ and $\tilde{d}_{1}$, it can be extended in a natural way from characteristic functions to any $B V$ function by using $d_{s t}$, and should provide a better approximation of the desired circuit than $d_{2}$, which seems to be too weak for this purpose.

### 3.4 Convolutions of characteristic functions and Gaussian distributions

We recall that $G(x)=(2 \pi)^{-1} \mathrm{e}^{-|x|^{2} / 2}, x \in \mathbb{R}^{2}$, and let us note that $\hat{G}(\xi)=\mathrm{e}^{-|\xi|^{2} / 2}, \xi \in \mathbb{R}^{2}$. Moreover, $\|G\|_{L^{1}\left(\mathbb{R}^{2}\right)}=1$. For any positive constant $s$ we denote by $G_{s}(x)=s^{-2} G(x / s)$, $x \in \mathbb{R}^{2}$. We note that $\left\|G_{s}\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}=1$ and $\widehat{G_{s}}(\xi)=\hat{G}(s \xi), \xi \in \mathbb{R}^{2}$.

Let $D$ be a bounded open set which is Lipschitz with constants $R_{0}$ and $L$ and let $\chi_{D}$ be its characteristic function. We investigate how $\chi_{D}$ is perturbed if we convolute it with $G$. We call $v=\chi_{D} * G$, that is

$$
v(x)=\int_{\mathbb{R}^{2}} \chi_{D}(x-y) G(y) \mathrm{d} y=\int_{\mathbb{R}^{2}} \chi_{D}(y) G(x-y) \mathrm{d} y, \quad x \in \mathbb{R}^{2}
$$

We recall that we shall use the positive constants $\delta_{1}, \delta_{2}$ and $m_{1}$, and the sets $M_{\delta_{1}}$ and $N_{\delta}$ introduced at the beginning of Subsection 3.3.

Proposition 3.7 Under the previous notation and assumptions, let us fix $\delta, 0<\delta \leq \delta_{2} / 4$. Then there exist constants $R_{0} \geq 1, \tilde{h}, 0<\tilde{h} \leq 1 / 24$, and $a_{1}>0$, depending on $L$ and $\delta$ only, such that the following estimates hold. For any $x \in \mathbb{R}^{2}$,

$$
\begin{equation*}
\text { if } \tilde{h}<v(x)<1-\tilde{h}, \text { then } x \in \bar{B}_{m_{1} \delta R_{0}}(\partial D), \tag{3.22}
\end{equation*}
$$

and for any $x \in \partial D$,

$$
\begin{equation*}
\text { if } y \in N_{\delta}(x) \text {, then } \nabla v(y) \cdot\left(-e_{2}\right) \geq a_{1} . \tag{3.23}
\end{equation*}
$$

Proof. If $x \notin \bar{B}_{m_{1} \delta R_{0}}(\partial D)$, then we have

$$
v(x) \leq \mathrm{e}^{-m_{1}^{2} \delta^{2} R_{0}^{2} / 2}, \text { if } x \notin D,
$$

and

$$
v(x) \geq 1-\mathrm{e}^{-m_{1}^{2} \delta^{2} R_{0}^{2} / 2}, \text { if } x \in D .
$$

Consequently, provided $\tilde{h}=\mathrm{e}^{-m_{1}^{2} \delta^{2} R_{0}^{2} / 2} \leq 1 / 24$, we may conclude that 3.22 holds.
Let us take $x \in \partial D$ and $y \in N_{\delta}(x)$. Then, denoting by $\nu$ the exterior unit normal vector to $D$,

$$
\nabla v(y) \cdot\left(-e_{2}\right)=\int_{\partial D} G(y-z) \nu(z) \cdot e_{2} \mathrm{~d} \mathcal{H}^{1}(z)
$$

Therefore,

$$
\begin{aligned}
& \nabla v(y) \cdot\left(-e_{2}\right)= \int_{\partial D \cap \bar{B}_{2 \delta R_{0}}(y)} G(y-z) \nu(z) \cdot e_{2} \mathrm{~d} \mathcal{H}^{1}(z)+ \\
& \int_{\partial D \cap\left(B_{R_{0} / 2}(y) \backslash \bar{B}_{2 \delta R_{0}}(y)\right)} G(y-z) \nu(z) \cdot e_{2} \mathrm{~d} \mathcal{H}^{1}(z)+ \\
& \quad \int_{\partial D \backslash B_{R_{0} / 2}(y)} G(y-z) \nu(z) \cdot e_{2} \mathrm{~d} \mathcal{H}^{1}(z)=A+B+C .
\end{aligned}
$$

Since $B_{R_{0} / 2}(y)$ is contained in $B_{R_{0}}(x)$, for any $z \in \partial D \cap B_{R_{0} / 2}(y)$, we have $\nu(z) \cdot e_{2} \geq$ $c_{1}>0$ where $c_{1}$ is a constant depending on $L$ only. Moreover, the length of $\partial D \cap \bar{B}_{2 \delta R_{0}}(y)$ is also bounded from below by $c_{2} \delta R_{0}, c_{2}>0$ depending on $L$ only. Therefore, we obtain that $A \geq c_{1} c_{2} \delta R_{0} \mathrm{e}^{-2 \delta^{2} R_{0}^{2}}$ and $B \geq 0$.

For what concerns the term $C$, with the help of (3.14), we can find a constant $C_{1}$, depending on $L$ only, such that, for any $R_{0} \geq 1$, we have

$$
|C| \leq C_{1} R_{0} \mathrm{e}^{-R_{0}^{2} / 8} .
$$

Therefore, we can find $R_{0} \geq 1$, depending on $L$ and $\delta$ only, such that $\tilde{h}=\mathrm{e}^{-m_{1}^{2} \delta^{2} R_{0}^{2} / 2} \leq$ $1 / 24$, and $2 C_{1} \mathrm{e}^{-R_{0}^{2} / 8} \leq c_{1} c_{2} \delta \mathrm{e}^{-2 \delta^{2} R_{0}^{2}}$. We set $a_{1}=(1 / 2) c_{1} c_{2} \delta R_{0} \mathrm{e}^{-2 \delta^{2} R_{0}^{2}}$ and the proof is concluded.

Remark 3.8 Without loss of generality, we may choose $R_{0}$ such that it also satisfies

$$
\begin{equation*}
\|\nabla G\|_{L^{1}\left(\mathbb{R}^{2} \backslash B_{R_{0} / 2}\right.} \leq(1 / 12) a_{1} \tag{3.24}
\end{equation*}
$$

In the sequel, we shall fix $\delta=\delta_{2} / 4$ and $R_{0}$ as the corresponding constant in Proposition 3.7 such that (3.24) holds. We note that, in this case, $\delta$ and $R_{0}$ depend on $L$ only. We shall also fix a constant $R \geq 10 R_{0}$. We recall that, with a slight abuse of notation, we identify $L^{1}\left(B_{R}\right)$ with the set of real valued $L^{1}\left(\mathbb{R}^{2}\right)$ functions that are equal to zero almost everywhere outside $B_{R}$. The same proof of Proposition 3.7 allows us to prove this corollary.

Corollary 3.9 For any $s, 0<s \leq 1$, let $r=s R_{0}$ and let $D$ be a bounded open set which is Lipschitz with constants $r$ and L. Let $v_{s}=\chi_{D} * G_{s}$. Then, for any $x \in \mathbb{R}^{2}$,

$$
\text { if } \tilde{h}<v_{s}(x)<1-\tilde{h}, \text { then } x \in \bar{B}_{m_{1} \delta r}(\partial D),
$$

and for any $x \in \partial D$,

$$
\text { if } y \in N_{\delta}(x) \text {, then } \nabla v_{s}(y) \cdot\left(-e_{2}\right) \geq a_{1} R_{0} / r=a_{1} / s
$$

We conclude this part with the following perturbation argument. Let us consider a function $\psi$ such that either $\psi \in C^{1}\left(\mathbb{R}^{2}\right) \cap W^{1,1}\left(\mathbb{R}^{2}\right)$ or $\psi \in W^{2,1}\left(\mathbb{R}^{2}\right)$ and that, for some $\tilde{\delta}>0$,

$$
\|\psi\|_{W^{1,1}\left(\mathbb{R}^{2}\right)} \leq \tilde{\delta} .
$$

Let $\tilde{G}=G+\psi$. Then the following result holds.
Corollary 3.10 Let us assume that $\tilde{\delta} \leq \min \left\{\tilde{h}, a_{1} / 2\right\}$.
For any $s, 0<s \leq 1$, let $r=s R_{0}$ and let $D$ be a bounded open set which is Lipschitz with constants $r$ and $L$. Let $v_{s}=\chi_{D} * \tilde{G}_{s}$. Then, for any $x \in \mathbb{R}^{2}$,

$$
\begin{equation*}
\text { if } 2 \tilde{h}<v_{s}(x)<1-2 \tilde{h}, \text { then } x \in \bar{B}_{m_{1} \delta r}(\partial D), \tag{3.25}
\end{equation*}
$$

and for any $x \in \partial D$,

$$
\begin{equation*}
\text { if } y \in N_{\delta}(x) \text {, then } \nabla v_{s}(y) \cdot\left(-e_{2}\right) \geq a_{1} R_{0} / 2 r=a_{1} /(2 s) \tag{3.26}
\end{equation*}
$$

Proof. It follows immediately from the previous corollary and Proposition 3.5. We first notice that in either cases $\chi_{D} * \tilde{G}_{s} \in C^{1}\left(\mathbb{R}^{2}\right)$. Moreover we have

$$
\left\|\chi_{D} * \tilde{G}_{s}-\chi_{D} * G_{s}\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq\left\|\chi_{D}\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}\left\|\psi_{s}\right\|_{L^{1}} \leq\|\psi\|_{L^{1}\left(\mathbb{R}^{2}\right)}
$$

and

$$
\begin{aligned}
&\left\|\nabla\left(\chi_{D} * \tilde{G}_{s}-\chi_{D} * G_{s}\right)\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}=\left\|\chi_{D} *\left(\nabla \psi_{s}\right)\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq \\
&\left\|\chi_{D}\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}\left\|\nabla \psi_{s}\right\|_{L^{1}} \leq\|\nabla \psi\|_{L^{1}\left(\mathbb{R}^{2}\right)} / s .
\end{aligned}
$$

Thus the conclusion follows.

## 4 Relationship between a mask and its image intensity

In this section we study the relationship between a function representing a mask (not necessarily a characteristic function of a domain) and its associated image intensity. We recall the notation used. We fix $\delta=\delta_{2} / 4$ and $R_{0}$ as the corresponding constant in Proposition 3.7 such that (3.24) holds. We note that, in this case, $\delta$ and $R_{0}$ depend on $L$ only. We shall also fix a constant $R \geq 10 R_{0}$. We recall that, with a slight abuse of notation, we identify $L^{1}\left(B_{R}\right)$ with the set of real valued $L^{1}\left(\mathbb{R}^{2}\right)$ functions that are equal to zero almost everywhere outside $B_{R}$. We recall that $A=\left\{u \in L^{1}\left(B_{R}\right): 0 \leq u \leq 1\right.$ a.e. in $\left.B_{R}\right\}$.

Fixed $\tilde{\delta}>0$, we assume that $\psi \in W^{2,1}\left(\mathbb{R}^{2}\right)$ and that

$$
\|\psi\|_{W^{1,1}\left(\mathbb{R}^{2}\right)} \leq \tilde{\delta}
$$

We denote $\tilde{G}=G+\psi$ and, for any $s, 0<s \leq 1$, we define the operator $\mathcal{T}_{s}: L^{1}\left(B_{R}\right) \rightarrow$ $W^{2,1}\left(\mathbb{R}^{2}\right)$ as follows

$$
\mathcal{T}_{s}(u)=u * \tilde{G}_{s}, \quad \text { for any } u \in L^{1}\left(B_{R}\right) .
$$

The point spread function we use, $T$, can be described in general by the function $\tilde{G}$. Therefore a study of properties of convolutions with $\tilde{G}$ will be useful.

We remark that the following continuity properties of the operator $\mathcal{T}_{s}$ hold. For any $p$, $1 \leq p \leq+\infty$, and any $u \in L^{1}\left(B_{R}\right)$, we have, for an absolute constant $C$,

$$
\begin{align*}
& \left\|\mathcal{T}_{s}(u)\right\|_{L^{p}\left(\mathbb{R}^{2}\right)} \leq\|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|u\|_{L^{p}\left(\mathbb{R}^{2}\right)},  \tag{4.1}\\
& \left\|\nabla \mathcal{T}_{s}(u)\right\|_{L^{p}\left(\mathbb{R}^{2}\right)} \leq(C / s)\|\nabla \tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|u\|_{L^{p}\left(\mathbb{R}^{2}\right)},  \tag{4.2}\\
& \left\|D^{2} \mathcal{T}_{s}(u)\right\|_{L^{p}\left(\mathbb{R}^{2}\right)} \leq\left(C / s^{2}\right)\left\|D^{2} \tilde{G}\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|u\|_{L^{p}\left(\mathbb{R}^{2}\right)} \tag{4.3}
\end{align*}
$$

Let $J \in C^{0}\left(\mathbb{R}^{2}\right)$. For any $u \in L^{1}\left(B_{R}\right)$, we define $U \in L^{1}\left(\mathbb{R}^{4}\right)$ as follows

$$
U(x, y)=u(x) u(y) J(x-y), \quad \text { for any } x, y \in \mathbb{R}^{2} .
$$

Then, for any $s, 0<s \leq 1$, we define $H_{s} \in W^{2,1}\left(\mathbb{R}^{4}\right)$ in the following way

$$
H_{s}(x, y)=\tilde{G}_{s}(x) \tilde{G}_{s}(y), \quad \text { for any } x, y \in \mathbb{R}^{2}
$$

Therefore, for any $p, 1 \leq p \leq+\infty$, and any $u \in L^{1}\left(B_{R}\right)$, we have, for an absolute constant C,

$$
\begin{align*}
& \left\|U * H_{s}\right\|_{L^{p}\left(\mathbb{R}^{4}\right)} \leq\|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}^{2}\|J\|_{L^{\infty}\left(B_{2 R}\right)}\|u\|_{L^{p}\left(\mathbb{R}^{2}\right)}^{2},  \tag{4.4}\\
& \left\|\nabla\left(U * H_{s}\right)\right\|_{L^{p}\left(\mathbb{R}^{4}\right)} \leq(C / s)\|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|\nabla \tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|J\|_{L^{\infty}\left(B_{2 R}\right)}\|u\|_{L^{p}\left(\mathbb{R}^{2}\right)}^{2},  \tag{4.5}\\
& \left\|D^{2}\left(U * H_{s}\right)\right\|_{L^{p}\left(\mathbb{R}^{4}\right)} \leq\left(C / s^{2}\right)\|\tilde{G}\|_{W^{2,1}\left(\mathbb{R}^{2}\right)}^{2}\|J\|_{L^{\infty}\left(B_{2 R}\right)}\|u\|_{L^{p}\left(\mathbb{R}^{2}\right)}^{2} . \tag{4.6}
\end{align*}
$$

Let us fix $p>4$ and let $\alpha=1-4 / p, 0<\alpha<1$. Then, if $u \in A$ we have $U * H_{s} \in C^{1, \alpha}\left(\mathbb{R}^{4}\right)$ and, for some absolute constant $C$ depending on $p$,

$$
\left\|U * H_{s}\right\|_{C^{1, \alpha}\left(\mathbb{R}^{4}\right)} \leq C\left\|U * H_{s}\right\|_{W^{2, p}\left(\mathbb{R}^{4}\right)}
$$

We define $\mathcal{P}_{J, s}: A \rightarrow C^{1, \alpha}\left(\mathbb{R}^{2}\right)$ and $\mathcal{P}_{1, s}: A \rightarrow C^{1, \alpha}\left(\mathbb{R}^{2}\right)$ as follows. For any $u \in A$

$$
\mathcal{P}_{J, s}(u)(x)=\int_{\mathbb{R}^{2}} \int_{\mathbb{R}^{2}} u(\xi) \tilde{G}_{s}(x-\xi) J(\xi-\eta) \tilde{G}_{s}(x-\eta) u(\eta) \mathrm{d} \xi \mathrm{~d} \eta, \quad x \in \mathbb{R}^{2}
$$

and

$$
\mathcal{P}_{1, s}(u)=\left(\mathcal{T}_{s}(u)\right)^{2} .
$$

We notice that the two definitions are consistent when $J \equiv 1$ and that

$$
\mathcal{P}_{J, s}(u)(x)=\left(U * H_{s}\right)(x, x), \quad x \in \mathbb{R}^{2}
$$

Putting together the previous estimates we obtain the following result. We recall that we have fixed a number $p>4$ and that $\alpha=1-4 / p$.

Proposition 4.1 Under the previous notation and assumptions, let $\varepsilon=\|J-1\|_{L^{\infty}\left(B_{2 R}\right)}$. Then for any $u \in A$ and any $s, 0<s \leq 1$, we have, for some absolute constant $C$ depending on $p$,

$$
\begin{aligned}
& \left\|\mathcal{P}_{J, s}(u)\right\|_{C^{0, \alpha}\left(\mathbb{R}^{2}\right)} \leq((1+\varepsilon) C / s)\|\tilde{G}\|_{W^{1,1}\left(\mathbb{R}^{2}\right)}^{2}\|u\|_{L^{p}\left(B_{R}\right)}^{2} \leq \\
& \quad((1+\varepsilon) C / s)\|\tilde{G}\|_{W^{1,1}\left(\mathbb{R}^{2}\right)}^{2}\|u\|_{L^{1}\left(B_{R}\right)}^{2 / p} .
\end{aligned}
$$

The same estimate holds also for the gradient, namely

$$
\begin{aligned}
&\left\|\nabla \mathcal{P}_{J, s}(u)\right\|_{C^{0, \alpha}\left(\mathbb{R}^{2}\right)} \leq\left((1+\varepsilon) C / s^{2}\right)\|\tilde{G}\|_{W^{2,1}\left(\mathbb{R}^{2}\right)}^{2}\|u\|_{L^{p}\left(B_{R}\right)}^{2} \leq \\
&\left.\quad(1+\varepsilon) C / s^{2}\right)\|\tilde{G}\|_{W^{2,1}\left(\mathbb{R}^{2}\right)}^{2}\|u\|_{L^{1}\left(B_{R}\right)}^{2 / p} .
\end{aligned}
$$

Furthermore, we have

$$
\begin{equation*}
\left\|\mathcal{P}_{J, s}(u)-\mathcal{P}_{1, s}(u)\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}=\left\|\mathcal{P}_{J-1, s}(u)\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq\|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}^{2} \varepsilon \tag{4.7}
\end{equation*}
$$

and, for some absolute constant $C$,

$$
\begin{equation*}
\left\|\nabla\left(\mathcal{P}_{J, s}(u)-\mathcal{P}_{1, s}(u)\right)\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq C\|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|\nabla \tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)} \varepsilon / s \tag{4.8}
\end{equation*}
$$

Although $\mathcal{P}_{J, s}$ is nonlinear in its argument $u$, by a simple adaptation of the previous reasonings, we obtain that for any $u_{1}, u_{2} \in A$, and for some absolute constant $C$ depending on $p$, we have the following corresponding estimates

$$
\begin{aligned}
\left\|\mathcal{P}_{J, s}\left(u_{1}\right)-\mathcal{P}_{J, s}\left(u_{2}\right)\right\|_{C^{0, \alpha}\left(\mathbb{R}^{2}\right)} \leq((1+\varepsilon) C / s)\|\tilde{G}\|_{W^{1,1}\left(\mathbb{R}^{2}\right)}^{2} R^{2 / p}\left\|u_{1}-u_{2}\right\|_{L^{p}\left(B_{R}\right)} \leq \\
2((1+\varepsilon) C / s)\|\tilde{G}\|_{W^{1,1}\left(\mathbb{R}^{2}\right)}^{2} R^{2 / p}\left\|u_{1}-u_{2}\right\|_{L^{1}\left(B_{R}\right)}^{1 / p}
\end{aligned}
$$

and

$$
\begin{array}{r}
\left\|\nabla\left(\mathcal{P}_{J, s}\left(u_{1}\right)-\mathcal{P}_{J, s}\left(u_{2}\right)\right)\right\|_{C^{0, \alpha}\left(\mathbb{R}^{2}\right)} \leq\left((1+\varepsilon) C / s^{2}\right)\|\tilde{G}\|_{W^{2,1}\left(\mathbb{R}^{2}\right)}^{2} R^{2 / p}\left\|u_{1}-u_{2}\right\|_{L^{p}\left(B_{R}\right)} \leq \\
2\left((1+\varepsilon) C / s^{2}\right)\|\tilde{G}\|_{W^{2,1}\left(\mathbb{R}^{2}\right)}^{2} R^{2 / p}\left\|u_{1}-u_{2}\right\|_{L^{1}\left(B_{R}\right)}^{1 / p} .
\end{array}
$$

Therefore, $\mathcal{P}_{J, s}: A \rightarrow C^{1, \alpha}\left(\mathbb{R}^{2}\right)$ is Lipschitz continuous with respect to the $L^{p}$ norm and Hölder continuous with exponent $1 / p$ with respect to the $L^{1}$ norm.

We fix $\tilde{\delta}$ such that $0<\tilde{\delta} \leq \min \left\{\tilde{h}, a_{1} / 2\right\}$, with $\tilde{h}, 0<\tilde{h} \leq 1 / 24$, and $a_{1}>0$ as in Proposition 3.7, thus depending on $L$ only. We define the corresponding $T$ and $s_{0}$ as in Lemma 2.2. We finally fix $s=(k \mathrm{NA})^{-1}, 0<s \leq 1 / s_{0}$, and $\sigma, 0<\sigma \leq s$. Then we define

$$
I_{s, \sigma}(u)=\mathcal{P}_{J, s s_{0}}(u), \quad \text { for any } u \in A,
$$

where $J$ is given by (2.3). We recall that in (2.2) we defined $K=T_{s s_{0}}$, therefore for any open set $D \subset B_{R}$ we have that $I_{s, \sigma}\left(\chi_{D}\right)$ is the light intensity on the image plane corresponding to the mask $D$, see (2.1).

We denote by $\mathcal{H}: \mathbb{R} \rightarrow \mathbb{R}$ the Heaviside function such that $\mathcal{H}(t)=0$ for any $t \leq 0$ and $\mathcal{H}(t)=1$ for any $t>0$. For any constant $h$ we set $\mathcal{H}_{h}(t)=\mathcal{H}(t-h)$ for any $t \in \mathbb{R}$. Then, for any $h, 0<h<1$, any $s, 0<s \leq 1 / s_{0}$, and any $\sigma, 0<\sigma \leq s$, we define the operator $\mathcal{W}: A \rightarrow L^{\infty}\left(\mathbb{R}^{2}\right)$ as follows

$$
\begin{equation*}
\mathcal{W}(u)=\mathcal{H}_{h}\left(I_{s, \sigma}(u)\right), \quad \text { for any } u \in A \tag{4.9}
\end{equation*}
$$

Clearly, for any $u \in A, \mathcal{W}(u)$ is the characteristic function of an open set, which we shall call $\Omega(u)$. That is

$$
\begin{equation*}
\Omega(u)=\left\{x \in \mathbb{R}^{2}: I_{s, \sigma}(u)(x)>h\right\}, \quad \text { for any } u \in A . \tag{4.10}
\end{equation*}
$$

In other words, $\chi_{\Omega(u)}=\mathcal{W}(u)=\mathcal{H}_{h}\left(I_{s, \sigma}(u)\right)$. Moreover, whenever $u=\chi_{D}$, where $D$ is an open set contained in $B_{R}$, we shall denote $\Omega(D)=\Omega\left(\chi_{D}\right)$.

The final, and crucial, result of this section is the following.

Theorem 4.2 Let us fix a positive constant L. Let $\delta=\delta_{2} / 4$ and let $R_{0}$ be as in Proposition 3.7 and such that (3.24) holds. Let us also fix $R \geq 10 R_{0}$ and $p, p>4$, and $\alpha=1-4 / p$.

We fix $\tilde{\delta}$ such that $0<\tilde{\delta} \leq \min \left\{\tilde{h}, a_{1} / 2\right\}$, with $\tilde{h}, 0<\tilde{h} \leq 1 / 24$, and $a_{1}>0$ as in Proposition 3.7, thus depending on $L$ only. We define the corresponding $T$ and $s_{0}$ as in Lemma 2.2. We finally fix $s=(k N A)^{-1}, 0<s \leq 1 / s_{0}$, and $\sigma, 0<\sigma \leq s$. Then, for any $u \in A$ we define

$$
I_{s, \sigma}(u)=\mathcal{P}_{J, s s_{0}}(u)
$$

where $J$ is given by (2.3).
Then for any $h, 1 / 3 \leq h \leq 2 / 3$, and any $s, 0<s \leq 1 / s_{0}$, we can find positive constants $\tilde{\sigma}_{0}, 0<\tilde{\sigma}_{0} \leq 1$, and $\gamma_{0}$, depending on $L, R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}$, $p$, and ss ${ }_{0}$ only, such that for any $\sigma, 0<\sigma \leq \tilde{\sigma}_{0} s$, and any $\gamma, 0<\gamma \leq \gamma_{0}$, the following holds.

Let $\mathcal{A}=\mathcal{A}^{0,1}(r, L, R)$, where $r=s s_{0} R_{0}$. Let $\tilde{R}=R+2 m_{1} \delta R_{0}$, where $m_{1}, 0<$ $m_{1} \leq 1$, depends on $L$ only. Then, for any $u \in \mathcal{A}_{\gamma}$, we have that $\Omega(u) \Subset B_{\tilde{R}}$ and $\Omega(u) \in \mathcal{A}^{1, \alpha}\left(r_{1}, L_{1}, \tilde{R}\right)$. Here $r_{1}=s s_{0} \tilde{R}_{0} \leq r$, where $\tilde{R}_{0} \leq \delta_{1} R_{0} / 8$ depends on $L$ only, whereas $L_{1} \geq L$ depends on $L, R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}, p$ and $s s_{0}$ only.

Moreover, the map $\mathcal{W}: \mathcal{A}_{\gamma} \rightarrow B V\left(B_{\tilde{R}}\right)$ is uniformly continuous with respect to the $L^{1}$ norm on $\mathcal{A}_{\gamma}$ and the distance $d_{\text {st }}$ on $B V\left(B_{\tilde{R}}\right)$.

Remark 4.3 We observe that the distance $d_{s t}$ in $B V\left(B_{\tilde{R}}\right)$ between $\mathcal{W}\left(u_{1}\right)$ and $\mathcal{W}\left(u_{2}\right)$ corresponds to the distance $d_{3}$ related to $B_{\tilde{R}}$ between $\Omega\left(u_{1}\right)$ and $\Omega\left(u_{2}\right)$.

Proof of Theorem 4.2. The proof is a consequence of the previous analysis. We fix $s$, $0<s \leq 1 / s_{0}$, and $h, 1 / 3 \leq h \leq 2 / 3$.

Let us begin with the following preliminary case. Let $u=\chi_{D}$, where $D \in \mathcal{A}$, and let $v=\mathcal{T}_{s s_{0}}(u)$ and $\tilde{W}=\mathcal{H}_{h}\left(\left(\mathcal{T}_{s s_{0}}(u)\right)^{2}\right)$. We apply Corollary 3.10 and we obtain the following results.

If $\tilde{\Omega}$ is the open set such that $\tilde{W}=\chi_{\tilde{\Omega}}$, then, by 3.25 , we notice that $\partial \tilde{\Omega} \subset \bar{B}_{m_{1} \delta r}(\partial D)$ and that $\left(D \backslash \bar{B}_{m_{1} \delta r}(\partial D)\right) \subset \tilde{\Omega}$ and $\left(\mathbb{R}^{2} \backslash \bar{B}_{m_{1} \delta r}(\bar{D})\right) \cap \tilde{\Omega}=\emptyset$. Therefore $\tilde{\Omega} \Subset B_{\tilde{R}}$.

We take any $x \in \partial D$ and any $y \in M_{\delta_{1} / 2}(x)$, with respect to the coordinate system depending on $x$. Then we consider the points $y^{-}=y-\delta r e_{2}$ and $y^{+}=y+\delta r e_{2}$. We have that $y^{ \pm} \in \partial N_{\delta}(x) \backslash \bar{B}_{m_{1} \delta r}(\partial D)$. Moreover, $y^{-} \in D$ and $v\left(y^{-}\right) \geq 11 / 12$, whereas $y^{+} \notin D$ and $v\left(y^{+}\right) \leq 1 / 12$. Let us call $\tilde{y}^{+}=y+t_{0} \delta r e_{2}$, where $t_{0} \in(-1,1], t_{0}$ depends on $y$, and $v\left(\tilde{y}^{+}\right)=1 / 12$ whereas $v\left(y+t \delta r e_{2}\right)<1 / 12$ for any $t \in\left(t_{0}, 1\right]$. Then we use (3.26) and we obtain that, for any $t \in\left[-1, t_{0}\right], v\left(y+t \delta r e_{2}\right) \geq 1 / 12$ and

$$
-\frac{\mathrm{d}}{\mathrm{~d} t}\left(v\left(y+t \delta r e_{2}\right)\right)^{2} \geq \delta r a_{1} /\left(12 s s_{0}\right)
$$

We may conclude that there exists a function $\varphi_{1}:\left[-\delta_{1} r / 2, \delta_{1} r / 2\right] \rightarrow \mathbb{R}$ such that, for any $y=\left(y_{1}, y_{2}\right) \in N_{\delta}(x)$ with $\left|y_{1}\right| \leq \delta_{1} r / 2,(v(y))^{2}=h$ if and only if $y_{2}=\varphi_{1}\left(y_{1}\right)$. We recall that

$$
\begin{equation*}
\|v\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq C_{1}, \quad\|\nabla v\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq C_{1} /\left(s s_{0}\right) \tag{4.11}
\end{equation*}
$$

where $C_{1}$ is an absolute constant, and

$$
\begin{equation*}
\|\nabla v\|_{C^{0, \alpha}\left(\mathbb{R}^{2}\right)} \leq C_{2} /\left(s s_{0}\right)^{2} \tag{4.12}
\end{equation*}
$$

where $C_{2}$ depends on $R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}$ and $p$ only.

We obtain that $v^{2}$ is a $C^{1, \alpha}$ function and, by the implicit function theorem, we conclude that the function $\varphi_{1}$ is actually $C^{1, \alpha}$. We observe that

$$
\left\|\varphi_{1}^{\prime}\right\|_{L^{\infty}\left[-\delta_{1} r / 2, \delta_{1} r / 2\right]} \leq C_{3},
$$

where $C_{3}$ depends on $L$ only. Without loss of generality, by a translation we may assume that $\varphi_{1}(0)=0$, thus $\left\|\varphi_{1}\right\|_{L^{\infty}\left[-\delta_{1} r / 2, \delta_{1} r / 2\right]} \leq C_{3} \delta_{1} r / 2$. Finally, for any $t_{1}, t_{2} \in\left[-\delta_{1} r / 2, \delta_{1} r / 2\right]$,

$$
\left|\varphi_{1}^{\prime}\left(t_{1}\right)-\varphi_{1}^{\prime}\left(t_{2}\right)\right| \leq\left(C_{4} /\left(s s_{0}\right)\right)\left|t_{1}-t_{2}\right|^{\alpha},
$$

where $C_{4}$ is a constant depending on $L, R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}$ and $p$ only.
Then, it is not difficult to prove that for some $r_{1}=s s_{0} \tilde{R}_{0} \leq r$, with $\tilde{R}_{0} \leq \delta_{1} R_{0} / 8$ depending on $L$ only, we can find $L_{1} \geq L$, depending on $L, R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}, p$ and $s s_{0}$ only, such that $\tilde{\Omega} \in \mathcal{A}^{1, \alpha}\left(r_{1}, L_{1}, \tilde{R}\right)$. Let us also remark that we have obtained that $\tilde{d}_{1}(\tilde{\Omega}, D) \leq \delta r$.

Let us call $\varepsilon=\varepsilon(s, \sigma)=\varepsilon(\sigma / s)=\|J-1\|_{L^{\infty}\left(B_{2 R}\right)}$. We notice that, as $\sigma / s \rightarrow 0^{+}$, we have that $\varepsilon$ goes to 0 as well. We also assume, without loss of generality, that $\varepsilon$ is increasing with respect to the variable $\sigma / s$. Let us recall that, for any $u \in A$, if $w=I_{s, \sigma}(u)$, with $0<s \leq 1 / s_{0}$ and $0<\sigma \leq s$, then

$$
\begin{equation*}
\|w\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq C_{5}, \quad\|\nabla w\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq C_{5} /\left(s s_{0}\right), \tag{4.13}
\end{equation*}
$$

and

$$
\begin{equation*}
\|\nabla w\|_{C^{0, \alpha}\left(\mathbb{R}^{2}\right)} \leq C_{6} /\left(s s_{0}\right)^{2} \tag{4.14}
\end{equation*}
$$

where $C_{5}$ is an absolute constant and $C_{6}$ depends on $R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}$ and $p$ only.
For positive constants $\tilde{\sigma}_{0}, 0<\tilde{\sigma}_{0} \leq 1$, and $\gamma_{0}$, to be precised later, let us fix $\sigma, 0<\sigma \leq$ $\tilde{\sigma}_{0} s$ and $\gamma, 0<\gamma \leq \gamma_{0}$. We take $u \in \mathcal{A}_{\gamma}, v=\mathcal{T}_{s s_{0}}(u), w=I_{s, \sigma}(u)$, and $D \in \mathcal{A}$ such that $\left\|u-\chi_{D}\right\|_{L^{1}\left(B_{R}\right)} \leq \gamma$. Then we use Proposition 4.1 to infer that

$$
\begin{aligned}
& \left\|I_{\sigma, s}(u)-\left(\mathcal{T}_{s s_{0}}\left(\chi_{D}\right)\right)^{2}\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq\left\|I_{s, \sigma}(u)-v^{2}\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)}+\left\|\left(\mathcal{T}_{s s_{0}}(u)\right)^{2}-\left(\mathcal{T}_{s s_{0}}\left(\chi_{D}\right)\right)^{2}\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq \\
& \|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}^{2} \varepsilon+\left(2 C /\left(s s_{0}\right)\right)\|\tilde{G}\|_{L^{1}\left(\mathbb{R}^{2}\right)}\|\tilde{G}\|_{W^{1,1}\left(\mathbb{R}^{2}\right)}\left\|u-\chi_{D}\right\|_{L^{p}\left(B_{R}\right)} \leq\left(C_{7} /\left(s s_{0}\right)\right)\left(\varepsilon+\gamma^{1 / p}\right),
\end{aligned}
$$

where $C$ is an absolute constant and consequently $C_{7}$ depends on $p$ only.
Analogously, we can prove that

$$
\left\|\nabla\left(I_{s, \sigma}(u)-\left(\mathcal{T}_{s s_{0}}\left(\chi_{D}\right)\right)^{2}\right)\right\|_{L^{\infty}\left(\mathbb{R}^{2}\right)} \leq\left(C_{8} /\left(s s_{0}\right)^{2}\right)\left(\varepsilon+\gamma^{1 / p}\right)
$$

where the constant $C_{8}$ depends on $\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}$ and $p$ only.
We now choose the positive constants $\tilde{\sigma}_{0}$ and $\gamma_{0}$ in such a way that

$$
2\left(C_{7} /\left(s s_{0}\right)\right)\left(\varepsilon\left(\tilde{\sigma}_{0}\right)+\gamma_{0}^{1 / p}\right) \leq 1 / 6
$$

and

$$
\left(C_{8} /\left(s s_{0}\right)\right)\left(\varepsilon\left(\tilde{\sigma}_{0}\right)+\gamma_{0}^{1 / p}\right) \leq a_{1} / 24
$$

Clearly, $\tilde{\sigma}_{0}$ and $\gamma_{0}$ depends on $L,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}, p$ and $s s_{0}$ only.

Then we can apply to $w=I_{s, \sigma}(u)$ and $\Omega=\Omega(u)$ the same analysis we have used for $v$ and $\tilde{\Omega}$ in the first part of this proof. We may therefore conclude that if $u \in \mathcal{A}_{\gamma}$ and $D \in \mathcal{A}$ is such that $\left\|u-\chi_{D}\right\|_{L^{1}\left(B_{R}\right)} \leq \gamma$, then $\Omega \Subset B_{\tilde{R}}, \tilde{d}_{1}(\Omega, D) \leq \delta r$ and, taken $r_{1}$ as before, possibly with a smaller $\tilde{R}_{0}$ still depending on $L$ only, we can find $L_{1} \geq L$, depending on $L$, $R,\left\|D^{2}(T-G)\right\|_{L^{1}\left(\mathbb{R}^{2}\right)}, p$ and $s s_{0}$ only, such that $\Omega \in \mathcal{A}^{1, \alpha}\left(r_{1}, L_{1}, \tilde{R}\right)$.

This kind of argument leads us also to show that $\Omega$ shares the same topological properties of $D$, that is for example $\Omega$ and $\partial \Omega$ have the same number of connected components of $D$ and $\partial D$, respectively.

It remains to show the uniform continuity property. We recall that the operator $\mathcal{P}_{J, s s_{0}}$ is Hölder continuous from $A$, with the $L^{1}\left(B_{R}\right)$ norm, into $C^{1, \alpha}\left(\mathbb{R}^{2}\right)$, with its usual norm. This means that there exists a constant $\tilde{C}$ such that for any $u_{1}$ and $u_{2} \in A$, if we call $w_{1}=I_{s, \sigma}\left(u_{1}\right)$ and $w_{2}=I_{s, \sigma}\left(u_{2}\right)$, then

$$
\left\|w_{1}-w_{2}\right\|_{C^{1, \alpha}\left(\mathbb{R}^{2}\right)} \leq \tilde{C}\left\|u_{1}-u_{2}\right\|_{L^{1}\left(B_{R}\right)}^{1 / p} .
$$

A simple application of the previous analysis allows us to prove this claim
Claim 1 There exists a function $g:[0,+\infty) \rightarrow[0,+\infty)$, which is continuous, increasing and such that $g(0)=0$, satisfying the following property. For any $u \in \mathcal{A}_{\gamma}$, for any $\varepsilon>0$ and any $x \in \mathbb{R}^{2}$ we have

$$
\begin{equation*}
\text { if } x \notin B_{\varepsilon}(\partial \Omega(u)) \text {, then }\left|I_{s, \sigma}(u)-h\right|>g(\varepsilon) \text {. } \tag{4.15}
\end{equation*}
$$

Let us now assume that $u_{1}$ and $u_{2}$ belong to $\mathcal{A}_{\gamma}$ and let us fix $\varepsilon>0$. We can find $\eta>0$ such that if $\left\|u_{1}-u_{2}\right\|_{L^{1}\left(B_{R}\right)} \leq \eta$, then $\left\|w_{1}-w_{2}\right\|_{L^{\infty}}\left(\mathbb{R}^{2}\right) \leq g(\varepsilon)$.

Let us now take $x \in \partial \Omega\left(u_{1}\right)$, that is $x \in \mathbb{R}^{2}$ such that $w_{1}(x)=h$. We infer that $\left|w_{2}(x)-h\right| \leq g(\varepsilon)$, therefore by the claim we deduce that $x \in B_{\varepsilon}\left(\partial \Omega\left(u_{2}\right)\right)$. That is $\partial \Omega\left(u_{1}\right) \subset$ $B_{\varepsilon}\left(\partial \Omega\left(u_{2}\right)\right)$. By symmetry, we conclude that $\tilde{d}_{1}\left(\Omega\left(u_{1}\right), \Omega\left(u_{2}\right)\right) \leq \varepsilon$. In other words, the map which to any $u \in \mathcal{A}_{\gamma}$ associates the open set $\Omega(u)$ is uniformly continuous with respect to the $L^{1}$ norm on $\mathcal{A}_{\gamma}$ and the distance $\tilde{d}_{1}$. However, we have shown in Subsection 3.3 that the distances $d_{1}, \tilde{d}_{1}, d_{2}$ and $d_{3}$ are topologically equivalent on $\mathcal{A}^{1, \alpha}\left(r_{1}, L_{1}, \tilde{R}\right)$, to which all $\Omega(u)$ belongs, for any $u \in \mathcal{A}_{\gamma}$. Therefore the map $\mathcal{A}_{\gamma} \ni u \rightarrow \Omega(u)$ is uniformly continuous with respect to the $L^{1}$ norm on $\mathcal{A}_{\gamma}$ and any of the distances $d_{1}, \tilde{d}_{1}, d_{2}$ and $d_{3}$ related to $B_{\tilde{R}}$.

We observe that

$$
d_{2}\left(\Omega\left(u_{1}\right), \Omega\left(u_{2}\right)\right)=\left\|\mathcal{W}\left(u_{1}\right)-\mathcal{W}\left(u_{2}\right)\right\|_{L^{1}\left(B_{\tilde{R}+1}\right)}=\left\|\mathcal{W}\left(u_{1}\right)-\mathcal{W}\left(u_{2}\right)\right\|_{L^{1}\left(B_{\vec{R}}\right)}
$$

whereas

$$
d_{3}\left(\Omega\left(u_{1}\right), \Omega\left(u_{2}\right)\right)=d_{2}\left(\Omega\left(u_{1}\right), \Omega\left(u_{2}\right)\right)+\left|P\left(\Omega\left(u_{1}\right)\right)-P\left(\Omega\left(u_{2}\right)\right)\right|=d_{s t}\left(\mathcal{W}\left(u_{1}\right), \mathcal{W}\left(u_{2}\right)\right)
$$

where $d_{s t}$ is here the distance inducing strict convergence in $B V\left(B_{\tilde{R}}\right)$. Therefore we conclude that $\mathcal{W}: \mathcal{A}_{\gamma} \rightarrow B V\left(B_{\tilde{R}}\right)$ is uniformly continuous with respect to the $L^{1}$ norm on $\mathcal{A}_{\gamma}$ and, on $B V\left(B_{\tilde{R}}\right)$, with respect either to the $L^{1}$ norm or to the $d_{s t}$ distance.

Remark 4.4 Let us finally remark that if, instead of taking $h \in[1 / 3,2 / 3]$, we simply assume $0<h<1$, then the same analysis may still be carried over. Clearly we need to change the values of $R_{0}$ and $\tilde{h}_{1}$ in Proposition 3.7, so that they depend on $h$ as well. As a consequence also the quantities introduced in the above Theorem 4.2 would depend on $h$.

## 5 Analysis of the inverse problem

Throughout this section, we shall keep the notation of Theorem4.2 and we shall also assume that the hypotheses of Theorem 4.2 are satisfied. We shall fix $h, 1 / 3 \leq h \leq 2 / 3$, and $s$, $0<s \leq 1 / s_{0}$, and we shall take $\sigma, 0<\sigma \leq \tilde{\sigma}_{0} s$, and $\gamma, 0<\gamma \leq \gamma_{0}, \tilde{\sigma}_{0}$ and $\gamma_{0}$ as in Theorem 4.2.

We call $\Omega_{0}$ the circuit to be reconstructed and we shall assume that it belongs to $\mathcal{A}=\mathcal{A}^{0,1}(r, L, R)$.

We recall that, by Proposition 3.6, $\mathcal{A}$ is compact with respect to the $d_{2}$ distance, which corresponds to the distance induced by the $L^{1}$ norm for the corresponding characteristic functions. Then it is an immediate consequence of the last part of Theorem 4.2, see also Remark 4.3 , that the problem

$$
\min _{D \in \mathcal{A}} d_{3}\left(\Omega(D), \Omega_{0}\right)
$$

admits a solution. We note that $\Omega(D)=\Omega\left(\chi_{D}\right)$ and that here $d_{3}$ is the distance defined in (3.10) related to $B_{\tilde{R}}$.

From a numerical point of view, the class $\mathcal{A}$ is rather difficult to handle. We try to reduce this difficulty by enlarging the class $\mathcal{A}$ to a class of characteristic functions of sets with finite perimeter. In order to keep the lower semicontinuity of the functional, we restrict ourselves to characteristic functions of sets with finite perimeter which are contained in $\mathcal{A}_{\gamma}$. Namely, we define the following functional $F_{0}: A \rightarrow[0,+\infty]$ such that for any $u \in A$ we have

$$
\begin{equation*}
F_{0}(u)=d_{s t}\left(\mathcal{W}(u), \chi_{\Omega_{0}}\right)+b P(u) \tag{5.1}
\end{equation*}
$$

where $P$ is the functional defined in (3.3) with $\mathcal{D}$ chosen to be $B_{R}, b$ is a positive parameter and $d_{s t}$ is the strict convergence distance in $B V\left(B_{\tilde{R}}\right)$. We recall that, whenever $u \in\{0,1\}$ almost everywhere in $B_{R}$ and $u \in B V\left(B_{R+1}\right)$, then $P(u)=P\left(u, B_{R+1}\right)=|D u|\left(B_{R+1}\right)$. Otherwise, $P(u)$, and consequently also $F_{0}(u)$, is equal to $+\infty$. Moreover, if $u \in \mathcal{A}_{\gamma}$, in particular if $u=\chi_{D}$ for some $D \in \mathcal{A}$, then $d_{s t}\left(\mathcal{W}(u), \chi_{\Omega_{0}}\right)=d_{3}\left(\Omega(u), \Omega_{0}\right)$, where again $d_{3}$ is the distance defined in (3.10) related to $B_{\tilde{R}}$.

We look for the solution to the following minimization problem

$$
\begin{equation*}
\min \left\{F_{0}(u): u \in \mathcal{A}_{\gamma}\right\} . \tag{5.2}
\end{equation*}
$$

We notice that such a minimization problem admits a solution.
Even if the class $\mathcal{A}_{\gamma}$ might still be not very satisfactory to handle from a numerical point of view, since it somehow involves handling the class $\mathcal{A}$, we believe that from a practical point of view such a restriction might be dropped and we might use the class $A \subset L^{1}\left(B_{R}\right)$ instead. In fact, we have a good initial guess, given by the target circuit $\chi_{\Omega_{0}}$, and it is reasonable to assume that the optimal mask will be a rather small perturbation of $\Omega_{0}$ itself. In fact, under our assumptions, by the arguments developed in the proof of Theorem 4.2, we can show that $\Omega(u)$ has the same topological properties of $D$, where $\chi_{D}$ is the element of $\mathcal{A}$ which is closest to $u$. Therefore if we look for a set $\Omega(u)$ as close as possible to $\Omega_{0}$, then at least we need to require that the set $D$ has the same topological properties of $\Omega_{0}$. For this reason and since $\Omega_{0} \in \mathcal{A}$, it might be essentially the same to perform the minimization in a small neighbourhood of $\mathcal{A}$ or in the whole $A$. On the other hand, again by our assumptions, we notice that whenever the boundary of $\Omega_{0}$ presents a corner, and this is often case, as
$\partial \Omega_{0}$ is often the union of a finite number of segments, then $\Omega_{0}$ cannot be reconstructed in an exact way, since $\Omega(u)$, for any $u \in \mathcal{A}_{\gamma}$, is a $C^{1, \alpha}$ set, thus its boundary cannot have any corner.

Besides dealing with the class $\mathcal{A}_{\gamma}$, there are several other difficulties. In particular, computing $F_{0}\left(\chi_{E}\right)$ for some $E \subset B_{R}$ is not an easy task, since it involves at least the computation of the perimeters of $E$ and of $\Omega\left(\chi_{E}\right)$. Furthermore, solving a minimization problem in the class of sets of finite perimeter is not a straightforward task from the numerical point of view.

In order to solve these difficulties, we use the following strategy. We approximate, in the sense of $\Gamma$-convergence, the functional $F_{0}$ with a family of functional $\left\{F_{\varepsilon}\right\}_{\varepsilon>0}$ which are easier to compute numerically and are defined on a set of smooth functions.

As in Section 2.3, we take a $C^{\infty}$ function $\phi: \mathbb{R} \rightarrow \mathbb{R}$ such that $\phi$ is nondecreasing, $\phi(t)=0$ for any $t \leq-1 / 2$ and $\phi(t)=1$ for any $t \geq 1 / 2$. For any $\eta>0$ and any $\tau \in \mathbb{R}$, let

$$
\phi_{\eta, \tau}(t)=\phi\left(\frac{t-\tau}{\eta}\right), \quad \text { for any } t \in \mathbb{R}
$$

Then we have the following result.
Proposition 5.1 For any $\eta>0$, let $\Phi_{\eta}: A \rightarrow C^{1, \alpha}\left(\mathbb{R}^{2}\right)$ be defined as

$$
\Phi_{\eta}(u)=\phi_{\eta, h}\left(I_{s, \sigma}(u)\right), \quad \text { for any } u \in A .
$$

Then, for any $\eta, 0<\eta \leq h, \Phi_{\eta}$ is Hölder continuous, with exponent $1 / p$, from $A$, with the $L^{1}\left(B_{R}\right)$ norm, into $C^{1, \alpha}\left(\mathbb{R}^{2}\right)$, with its usual norm.

Furthermore, as $\eta \rightarrow 0^{+}$, $\left(\mathcal{W}-\Phi_{\eta}\right): \mathcal{A}_{\gamma} \rightarrow B V\left(B_{\tilde{R}}\right)$ converges uniformly to zero on $\mathcal{A}_{\gamma}$ with respect to the distance $d_{\text {st }}$ on $B V\left(B_{\tilde{R}}\right)$.

Proof. The continuity property of $\Phi_{\eta}$ immediately follows by the continuity of $\mathcal{P}_{J, s s_{0}}$ and by the properties of $\phi$. We just note that the Hölder exponent is fixed, whereas the Hölder constant might depend upon $\eta$.

About the convergence result, we begin by recalling that $\mathcal{W}(u)=\mathcal{H}_{h}\left(I_{s, \sigma}(u)\right), u \in \mathcal{A}_{\gamma}$. We use Claim 1 introduced in the proof of Theorem 4.2. We call $t_{0}=\min (1, \sup \{g(s)$ : $s \in[0,+\infty)\}$ ) and $s_{0}$ the positive real number such that $g\left(s_{0}\right)=t_{0} / 2$. We call $g^{-1}$ : $\left[0, t_{0} / 2\right] \rightarrow\left[0, s_{0}\right]$ the continuous, increasing function which is the inverse of $g$ on such intervals. For any $\eta, 0<\eta \leq t_{0}$, we infer that $\Phi_{\eta}(u)(x)$ might be different from $\mathcal{W}(u)(x)$ only if $x \in B_{g^{-1}(\eta / 2)}(\partial \Omega(u))$. By estimates like (3.15) and 3.16), which are independent of $u \in \mathcal{A}_{\gamma}$, we obtain that $\left\|\left(\mathcal{W}-\Phi_{\eta}\right)(u)\right\|_{L^{1}\left(B_{\vec{R}}\right)}$ converges to zero, as $\eta \rightarrow 0^{+}$, uniformly for $u \in \mathcal{A} \gamma$.

For any $t \in \mathbb{R}$ and any $u \in \mathcal{A}_{\gamma}$, we call

$$
P(u, t)=P\left(\left\{x \in \mathbb{R}^{2}: I_{s, \sigma}(u)(x)>t\right\}, B_{\tilde{R}}\right) .
$$

It remains to prove that, as $\eta \rightarrow 0^{+},\left|D\left(\Phi_{\eta}(u)\right)\right|\left(B_{\tilde{R}}\right)=\int_{B_{\tilde{R}}}\left|\nabla\left(\Phi_{\eta}(u)\right)\right|$ converges to $|D(\mathcal{W}(u))|\left(B_{\tilde{R}}\right)=P(u, h)$ uniformly for $u \in \mathcal{A}_{\gamma}$. We argue in the following way. We have that, for any $\eta, 0<\eta \leq h$,

$$
\left|D\left(\Phi_{\eta}(u)\right)\right|\left(B_{\tilde{R}}\right)=\int_{B_{\tilde{R}}}\left|\nabla\left(\Phi_{\eta}(u)\right)\right|=\int_{B_{\tilde{R}}}\left|\phi_{\eta, h}^{\prime}\left(I_{s, \sigma}(u)\right)\right|\left|\nabla\left(I_{s, \sigma}(u)\right)\right| .
$$

Since $\phi_{\eta, h}^{\prime} \geq 0$, and for $\eta$ small enough, uniformly with respect to $u \in \mathcal{A}_{\gamma}, \phi_{\eta, h}^{\prime}\left(I_{s, \sigma}(u)\right)=0$ outside $B_{\tilde{R}}$, without loss of generality, we have that

$$
\left|D\left(\Phi_{\eta}(u)\right)\right|\left(B_{\tilde{R}}\right)=\int_{\mathbb{R}^{2}} \phi_{\eta, h}^{\prime}\left(I_{s, \sigma}(u)\right)\left|\nabla\left(I_{s, \sigma}(u)\right)\right| .
$$

By the coarea formula,

$$
\left|D\left(\Phi_{\eta}(u)\right)\right|\left(B_{\tilde{R}}\right)=\int_{-\infty}^{+\infty}\left(\int_{\left\{I_{s, \sigma}(u)=t\right\}} \phi_{\eta, h}^{\prime}(t) \mathrm{d} \mathcal{H}^{1}(y)\right) \mathrm{d} t .
$$

Therefore,

$$
\left|D\left(\Phi_{\eta}(u)\right)\right|\left(B_{\tilde{R}}\right)=\frac{1}{\eta} \int_{-\infty}^{+\infty} \phi^{\prime}\left(\frac{(t-h)}{\eta}\right) P(u, t) \mathrm{d} t=\int_{-1 / 2}^{+1 / 2} \phi^{\prime}(s) P(u, s \eta+h) \mathrm{d} s
$$

Since $|D(\mathcal{W}(u))|\left(B_{\tilde{R}}\right)=P(u, h)$ and $\int_{-1 / 2}^{+1 / 2} \phi^{\prime}(s) \mathrm{d} s=1$, we obtain that

$$
\left|\left|D\left(\Phi_{\eta}(u)\right)\right|\left(B_{\tilde{R}}\right)-|D(\mathcal{W}(u))|\left(B_{\tilde{R}}\right)\right| \leq \int_{-1 / 2}^{+1 / 2} \phi^{\prime}(s)|P(u, s \eta+h)-P(u, h)| \mathrm{d} s
$$

It remains to show that, as $\eta \rightarrow 0^{+}, \sup \{|P(u, t+h)-P(u, h)|: t \in[-\eta / 2,+\eta / 2]\}$ goes to zero uniformly with respect to $u \in \mathcal{A}_{\gamma}$. Therefore the proof is concluded by using the following claim.

Claim 2 There exist a positive constant $\eta_{0}$ and a continuous, increasing function $g_{1}$ : $\left[0, \eta_{0}\right] \rightarrow[0,+\infty)$, such that $g_{1}(0)=0$, such that for any $\eta, 0<\eta \leq \eta_{0}$, and any $u \in \mathcal{A}_{\gamma}$, we have that

$$
\sup \{|P(u, t+h)-P(u, h)|: t \in[-\eta / 2,+\eta / 2]\} \leq g_{1}(\eta)
$$

The proof of Claim 2 is a straightforward, although maybe lengthy to describe, consequence of the analysis developed in the proof of Theorem 4.2. We leave the details to the reader. We just notice that Claim 2 is a sort of generalization of Claim 1 and the arguments used to prove the two claims are essentially analogous.

We are now in the position of describing the approximating functionals and proving the $\Gamma$-convergence result. Let us a fix a constant $p_{1}, 1<p_{1}<+\infty$, and a continuous function $W: \mathbb{R} \rightarrow[0,+\infty)$ such that $W(t)=0$ if and only if $t \in\{0,1\}$. Let us denote by $P_{\varepsilon}, \varepsilon>0$, the functional defined in (3.2) with $p=p_{1}$, the function $W$ and $\mathcal{D}=B_{R}$. We recall that the functional $P$ is defined in (3.3), again with $\mathcal{D}=B_{R}$.

Then, for any $\varepsilon>0$, let us define $F_{\varepsilon}: A \rightarrow[0,+\infty]$ such that for any $u \in A$ we have

$$
\begin{equation*}
F_{\varepsilon}(u)=d_{s t}\left(\Phi_{\eta(\varepsilon)}(u), \chi_{\Omega_{0}}\right)+b P_{\varepsilon}(u) \tag{5.3}
\end{equation*}
$$

where $\eta:[0,+\infty) \rightarrow[0,+\infty)$ is a continuous, increasing function such that $\eta(0)=0$.
By the direct method, we can prove that each of the functionals $F_{\varepsilon}, \varepsilon>0$, admits a minimum either over $A$ or over $\mathcal{A}_{\gamma}$.

The $\Gamma$-convergence result is the following.

Theorem 5.2 Let us consider the metric space $(X, d)$ where $X=\mathcal{A}_{\gamma}$ and $d$ is the metric induced by the $L^{1}$ norm. Then, as $\varepsilon \rightarrow 0^{+}, F_{\varepsilon} \Gamma$-converges to $F_{0}$ on $X$ with respect to the distance $d$.

Proof. Let us fix a sequence $\left\{\varepsilon_{n}\right\}_{n=1}^{\infty}$ of positive numbers converging to zero as $n \rightarrow \infty$. Let, for any $n \in \mathbb{N}, F_{n}=F_{\varepsilon_{n}}$. We need to prove that $\Gamma-\lim _{n} F_{n}=F$. Let us also remark that we may extend $F_{n}$ and $F$ over $L^{1}\left(\mathbb{R}^{2}\right)$ by setting them equal to $+\infty$ outside $\mathcal{A}_{\gamma}$. Let us define $\tilde{P}_{\varepsilon}, \varepsilon>0$, and $\tilde{P}$ as the functionals which are equal to the functionals $P_{\varepsilon}$ and $P$, respectively, on $\mathcal{A}_{\gamma}$ and $+\infty$ elsewhere. We recall that $P_{\varepsilon}, \varepsilon>0$, and $P$ are defined in (3.2) and in (3.3), respectively, with $p=p_{1}$ and $\mathcal{D}=B_{R}$.

We observe that, as a consequence of Proposition 5.1 and of the stability of $\Gamma$-convergence under uniformly converging continuous perturbations, it is enough to show that $\Gamma$ - $\lim _{n} \tilde{P}_{n}=$ $\tilde{P}$, where $\tilde{P}_{n}=\tilde{P}_{\varepsilon_{n}}, n \in \mathbb{N}$. Let us prove this $\Gamma$-convergence result.

The $\Gamma$-liminf inequality is an immediate consequence of Theorem 3.1 and of the fact that $\mathcal{A}_{\gamma}$ is a closed subset of $L^{1}\left(B_{R}\right)$.

For what concerns the recovery sequence, then we argue in the following way. If $u \in A$ is such that $\left\|u-\chi_{D}\right\|_{L^{1}\left(B_{R}\right)}<\gamma$, for some $D \in \mathcal{A}$, then we again use Theorem 3.1 to construct a recovery sequence for such a function $u$, that is a sequence $\left\{u_{n}\right\}_{n=1}^{\infty}$ contained in $\mathcal{A}_{\gamma}$ such that, as $n \rightarrow \infty, u_{n} \rightarrow u$ in $L^{1}\left(B_{R}\right)$ and $\tilde{P}_{n}\left(u_{n}\right) \rightarrow \tilde{P}(u)$.

It remains to study the case when $u \in \partial \mathcal{A}_{\gamma}$ and $\tilde{P}(u)<+\infty$. In this case, we have that $u=\chi_{E}$, where $E \subset B_{R}$ is a set of finite perimeter, and we pick $D \in \mathcal{A}$ such that $\left\|\chi_{E}-\chi_{D}\right\|_{L^{1}\left(B_{R}\right)}=|E \Delta D|=\gamma$. Then at least one of these two cases must be satisfied. Either there exists $x \in B_{R} \backslash \bar{D}$ such that

$$
\lim _{\rho \rightarrow 0^{+}} \frac{\left|E \cap B_{\rho}(x)\right|}{\left|B_{\rho}(x)\right|}=1
$$

or there exists $x \in D$ such that

$$
\lim _{\rho \rightarrow 0^{+}} \frac{\left|E \cap B_{\rho}(x)\right|}{\left|B_{\rho}(x)\right|}=0 .
$$

We choose an arbitrary sequence $\left\{\rho_{j}\right\}_{j=1}^{\infty}$ of positive numbers such that $\lim _{j} \rho_{j}=0$. In the first case, for any $j \in \mathbb{N}$, we choose $E_{j}$ such that $\chi_{E_{j}}=\chi_{E}\left(1-\chi_{B_{\rho_{j}}(x)}\right)$. In the second case, we choose $E_{j}$ such that $\chi_{E_{j}}=\chi_{E}\left(1-\chi_{B_{\rho_{j}}(x)}\right)+\chi_{B_{\rho_{j}}(x)}$. We notice that, in either cases, for any $j \in \mathbb{N}, E_{j}$ is a set of finite perimeter such that $\left\|\chi_{E_{j}}-\chi_{D}\right\|_{L^{1}\left(B_{R}\right)}<\gamma$. Furthermore, as $j \rightarrow \infty$ we have that $\chi_{E_{j}} \rightarrow \chi_{E}$ in $L^{1}\left(B_{R}\right)$ and $P\left(E_{j}\right) \rightarrow P(E)$, that is $\tilde{P}\left(\chi_{E_{j}}\right) \rightarrow \tilde{P}\left(\chi_{E}\right)$. Then the proof may be concluded by following the arguments of Section 4.2 in [2] which we have briefly recalled in the proof of Theorem 3.1.

We remark that $\Omega_{0} \in \mathcal{A}$, therefore we may find a family $\left\{\tilde{u}_{\varepsilon}\right\}_{\varepsilon>0}$ such that, as $\varepsilon \rightarrow 0^{+}$, $\tilde{u}_{\varepsilon} \rightarrow \chi \Omega_{0}$ in $L^{1}\left(B_{R}\right)$ and $P_{\varepsilon}\left(\tilde{u}_{\varepsilon}\right) \rightarrow P\left(\Omega_{0}\right)$. Without loss of generality, we may assume that, for any $\varepsilon>0,0 \leq \tilde{u}_{\varepsilon} \leq 1$ almost everywhere in $B_{R}$ and that $\tilde{u}_{\varepsilon} \in \mathcal{A}_{\gamma}$. By Proposition 5.1, we conclude that $F_{\varepsilon}\left(\tilde{u}_{\varepsilon}\right) \rightarrow F_{0}\left(\Omega_{0}\right)<+\infty$. We obtain that for any $\varepsilon_{0}>0$ there exists a constant $C_{1}$ such that

$$
\begin{equation*}
\min _{\mathcal{A}_{\gamma}} F_{\varepsilon} \leq C_{1} \quad \text { for any } \varepsilon, 0<\varepsilon \leq \varepsilon_{0} \tag{5.4}
\end{equation*}
$$

Obviously, the same property is shared by the minimum values of $F_{\varepsilon}$ over $A$.
It remains to prove that the functionals $F_{\varepsilon}$ are equicoercive over $\mathcal{A}_{\gamma}$, that is that the following result holds.

Proposition 5.3 For any $\varepsilon_{0}>0$, there exists a compact subset $\mathcal{K}$ of $\mathcal{A}_{\gamma}$ such that for any $\varepsilon, 0<\varepsilon \leq \varepsilon_{0}$, we have

$$
\min _{\mathcal{K}} F_{\varepsilon}=\min _{\mathcal{A}_{\gamma}} F_{\varepsilon} .
$$

Proof. Let us take the constant $C_{1}$ as in (5.4). Let $u_{\varepsilon} \in \mathcal{A}_{\gamma}, 0<\varepsilon \leq \varepsilon_{0}$, be such that $F_{\varepsilon}\left(u_{\varepsilon}\right)=\min _{\mathcal{A}_{\gamma}} F_{\varepsilon}$. Then we observe that the set $\left\{u_{\varepsilon}\right\}_{0<\varepsilon \leq \varepsilon_{0}}$ satisfies the properties of Remark 3.4 for some constant $C$. Therefore $\left\{u_{\varepsilon}\right\}_{0<\varepsilon \leq \varepsilon_{0}}$ is precompact in $L^{1}\left(B_{R}\right)$ and the proof is concluded.

Remark 5.4 With an analogous proof, the same result of Proposition 5.3 holds if we replace $\mathcal{A}_{\gamma}$ with $A$.

By Theorem 5.2 and Proposition 5.3, we can apply the Fundamental Theorem of $\Gamma$ convergence to conclude with the following result.

Theorem 5.5 We have that $F_{0}$ admits a minimum over $\mathcal{A}_{\gamma}$ and

$$
\min _{\mathcal{A}_{\gamma}} F_{0}=\lim _{\varepsilon \rightarrow 0^{+}} \inf _{\mathcal{A}} F_{\varepsilon}=\lim _{\varepsilon \rightarrow 0^{+}} \min _{\mathcal{A}_{\gamma}} F_{\varepsilon} .
$$

Let $\varepsilon_{n}, n \in \mathbb{N}$, be a sequence of positive numbers converging to 0 . For any $n \in \mathbb{N}$, let $F_{n}=F_{\varepsilon_{n}}$. If $\left\{u_{n}\right\}_{n=1}^{\infty}$ is a sequence contained in $\mathcal{A}_{\gamma}$ which converges, as $n \rightarrow \infty$, to $u \in \mathcal{A}_{\gamma}$ in $L^{1}\left(B_{R}\right)$ and satisfies $\lim _{n} F_{n}\left(u_{n}\right)=\lim _{n} \inf _{\mathcal{A}_{\gamma}} F_{n}$, then $u$ is a minimizer for $F_{0}$ on $\mathcal{A}_{\gamma}$, that is $u$ solves the minimization problem (5.2).

We conclude with the following remark. With the notation of Theorem 5.5. if $\left\{u_{n}\right\}_{n=1}^{\infty}$ is a sequence contained in $\mathcal{A}_{\gamma}$ which satisfies $\lim _{n} F_{n}\left(u_{n}\right)=\lim _{n} \inf _{\mathcal{A}_{\gamma}} F_{n}$, then, by Remark 3.4 , we have that, up to passing to a subsequence, $\left\{u_{n}\right\}_{n=1}^{\infty}$ actually converges, as $n \rightarrow \infty$, to some function $u \in \mathcal{A}_{\gamma}$ in $L^{1}\left(B_{R}\right)$.

## 6 Discussion

We have provided a mathematical study of the inverse problem of photolithography. The approach we propose is to seek an approximate solution by formulating the geometrical problem using a phase-field method. We further relax the hard threshold involved in image exposure with an approximate Heaviside function. We show that the variational problem for the approximate solution is well-posed. This opens a way into designing mathematically rigorous numerical methods. We further show that as the approximation parameter goes to zero, a theoretical limit, the original optimization problem involving geometry is recovered.

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