## Article

# Identification of parameters in parabolic partial differential equation from final observations using deep learning 

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#### Abstract

In this work, we propose a deep learning approach for identifying parameters (initial condition, a coefficient in the diffusion term and source function) in parabolic partial differential equations (PDEs) from scattered final observations in space and noisy a priori knowledge. In Particular, we approximate the unknown solution and parameters by four deep neural networks trained to satisfy the differential operator, boundary conditions, a priori knowledge and observations. The proposed algorithm is mesh-free, which is key since meshes become infeasible in higher dimensions due to the number of grid points explosion. Instead of forming a mesh, the neural networks are trained on batches of randomly sampled time and space points. This work is devoted to the identification of several parameters of PDEs at the same time. The classical methods require a total a priori knowledge which is not feasible. While they cannot solve this inverse problem given such partial data, the deep learning method allows them to resolve it using minimal a priori knowledge.


Keywords: Deep learning; Heat equation; Hybrid method; Inverse problem; Model-driven solution; Neural networks; Optimization; Tikhonov regularization.

MSC: 35C05

## 1. Introduction

According to J.B. Keller [1], two problems are said to be inverses of each other if the formulation of one put the other involved. A more operational definition is that an inverse problem is to determine causes knowing effects. Thus, this problem is the inverse of the one called direct problem, consisting in deducing the effects, the causes being known.

In a large number of applications based on partial differential equations, we are interested in the identification of some properties of the system from partial knowledge. We consider the following boundary problem:

$$
\left\{\begin{array}{l}
\left.\partial_{t} u-\partial_{x}\left(a \partial_{x} u\right)=f, \forall x \in \Omega, \forall t \in\right] 0, T[,  \tag{1}\\
\left.\left.u\right|_{\partial \Omega}=g(t), \forall t \in\right] 0, T[ \\
u(0, x)=u_{0}(x), \forall x \in \Omega .
\end{array}\right.
$$

Heat equation (diffusion equation) is the archetype of parabolic models. In the unidimensional problem case, (1) models the evolution of the temperature $u(t, x)$ in a conductor bar, the initial temperature distribution at $t=0$ is given by the function $u_{0}(x)$, on the edge the temperature is maintained at $g(t)$, the heat sources are modeled by the given function $f(t, x)$, the specific heat $a(x)$ of a material is the amount of heat energy that it takes to raise one unit of mass of the material by one unit of temperature.

There are several inverse problems or identification problems associated with the heat equation. Given temperature $u$ (some measures of $u(t, x)$ ), each question among the following models an inverse problem:

- What are the boundary conditions which we must supply this system to have the temperature $u$ ?
- What is the initial condition to reach the temperature $u$ ?
- What is the source we should provide to the system to have $u$ as a solution of (1)?

Sometimes an inverse problem is modeled by a combination of questions above. We should notice that for other values of observations we can get other parameters. On other hand we can search $a(x)$, which we call an identification problem, since it is a characterization which identifies the material. Whatever observations we get the same coefficient value.

The reformulation of an inverse problem as minimization of an error between the real measures and the solution of the direct problem is the most important step. We seek to identify a parameter $P \in E_{P}$ which can be the initial condition $u_{0}$, the source $f$ or the coefficient $a$, from an observed function $u^{o b s}$ (measures). The following equation of state relates implicitly the parameter $P$ and the state $u$ :

$$
\begin{equation*}
F(P, u)=0 . \tag{2}
\end{equation*}
$$

Let

$$
\begin{equation*}
u=\phi(P), \tag{3}
\end{equation*}
$$

be the solution of the problem (2).
The inverse problem consists in identifying $P$ such as:

$$
\begin{equation*}
\phi(P)=u^{o b s} . \tag{4}
\end{equation*}
$$

The application $\phi$ is defined implicitly. It is non-linear, even though the equation of state (2) is linear. This obviously makes it more difficult to solve the inverse problem. We refer the interested reader in inverse problems to see $[2,3]$.

What we said previously suggests that equation (4) may not have a solution, and even if it does, the inverse application is not necessarily continuous. We are therefore going to introduce a formulation, a priori weaker, which has proved its usefulness. We replace equation (4) with the following minimization problem:

$$
\left\{\begin{array}{l}
\text { Find } P \in A_{a d} \text { such as } \\
J_{T}(P)=\min _{h \in A_{a d}} J_{T}(h),
\end{array}\right.
$$

where the space $A_{a d}$ is the admissible set of $P$ and the cost function $J_{T}$ is defined as follows:

$$
J_{T}(h)=\frac{1}{2 T}\left\|\phi(h)-u^{o b s}\right\|^{2} .
$$

This formulation is called a least squares method, and $J_{T}$ is the cost function. The observation being given once and for all, to evaluate the functional $J_{T}$ in one parameter $P$, we start by solving the equation of state (2), then the observation equation 4, and we compare the simulated observation to the measured one. Regarding the observation, measurements are made for a few points of space and at some points in time. So, we need to do an interpolation.

The first section is devoted to the presentation of classical methods of resolution, their limitations and our motivation. In the second section we prove that the minimization problem is well-posed. We will calculate the gradient of $J_{T}$ in the third section which we will use for classical methods in the fourth one. The fifth section is devoted to the presentation of our approach. We will prove, under certain conditions, the convergence of the $L^{2}$ error defined as (24) using the smoothness of the neural network approximations in the sixth section. We will present in the following section some numerical results of our algorithm which we will compare with classical algorithms in the last section.

## 2. Related work

The mathematical model leads to a non-convex minimization problem. The objective function can have several local minimums, as a result a descent algorithm ends at the meeting of the first local minimum. To solve this problem we can apply two different approaches.

The first approach is Tikhonov regularization which consists in defining the solution of the inverse problem as being the minimum of one least squares criterion penalized by a quadratic term:

$$
J_{T}(h)=\frac{1}{2 T}\left\|\phi(h)-u^{o b s}\right\|^{2}+\frac{\varepsilon}{2}\left\|h-h^{b}\right\|^{2},
$$

where $\varepsilon$ is the regularization coefficient and $h^{b}$ is a priori knowledge (background) of exact parameter $h$.

An essential difficulty with this approach stems from the fact that it requires a priori knowledge at all points of the domain of resolution. Moreover, in practice, the knowledge $h^{b}$ is partial (30\%), the size of $h^{b}$ is much less than that of $h$. This makes it difficult to find the regularization coefficient.

To solve this problem a second approach called hybrid method is proposed in [4,5] . The idea of this method is to marry the genetic algorithm with a gradient descent method. The individual in the genetic algorithm is the parameter $h$. Let's start with an initial population, we apply the genetic algorithm with its operators (crossing and mutation) to produce a new population of which each individual serves as a starting point for the steep descent method. This method is very effective for inverse problems associated with identifying a single parameter of PDEs for example $a$ or $u_{0}$ or $f$. For the identification of several parameters of the equation this method does not solve the problem of partial knowledge.

However, there are some drawbacks associated with those methods. The gradient of functional $J_{T}$ is computed by using adjoint method, which requires the resolution of another boundary problem whose parameters are dependent on $u$ (the solution of the main problem). The first drawback is that they require, in each descent iteration, resolution of two PDEs by classical methods (finite difference, finite element, .etc). It means that to solve the inverse problem we should solve direct problem several times and manage its complexity which makes those methods costly.

The second drawback is that they are not meshfree. Their application for high dimensional problems and complex domains is difficult since classical methods become infeasible in higher dimensions due to the explosion in the number of grid points. They also require using interpolation and numerical integrals.

The last drawback is that they require a complete a priori knowledge. The solution reconstruction cannot be made from partial data.

Recently, deep learning emerged as a powerful technique in many applications. It is used in many works related to PDEs like [6-8] where authors solve many types of PDEs by deep learning method. This success inspires us to apply such a technique to identifying parameters of PDEs from final observations.

This work is devoted to the identification of $a, u_{0}$ and $f$ of (1) from final observations by a new approach based on deep learning method which is meshfree and does not require solving direct problem, interpolation or numeric integrals. This new approach requires only partial a priori knowledge.

The sought after solution is driven by the model since the relation between all parameters of the system (diffusion phenomenon model) is the main information for the resolution.

## 3. Well-posedness

In this section we consider the case of homogeneous boundary data, i.e., $g(t)=0$.

We specify some notations we shall use.

$$
\begin{gathered}
S_{a}=S_{u_{0}}=L^{2}(\Omega), \quad S_{f}=L^{2}(] 0, T\left[; L^{2}(\Omega)\right) \\
S=S_{a} \times S_{u_{0}} \times S_{f}
\end{gathered}
$$

Let $P=\left(a, u_{0}, f\right) \in S$

$$
\|P\|_{S}=\|a\|_{S_{a}}+\left\|u_{0}\right\|_{S_{u_{0}}}+\|f\|_{S_{f}}
$$

Definition 1. (Variational formulation [9])
The variational formulation of the problem (1) is the following:
Find $u(t)$ function of $] 0, T\left[\right.$ in $H_{0}^{1}(\Omega)$ such as:

$$
\left\{\begin{array}{l}
\left.\frac{d}{d t}\langle u(t), v\rangle_{L^{2}(\Omega)}+\left\langle a u_{x}(t), v_{x}\right\rangle_{L^{2}(\Omega)}=\langle f(t), v\rangle_{L^{2}(\Omega)}, \forall v \in H_{0}^{1}(\Omega), \forall t \in\right] 0, T[,  \tag{5}\\
u(x, 0)=u_{0}(x), \forall x \in \Omega .
\end{array}\right.
$$

Definition 2. (Weak solution)
A weak solution of (1) is a solution of the variational formulation 5 .
Let $a \in L^{\infty}(\Omega)$ where $a(x)>\gamma, \gamma>0, \forall x \in \bar{\Omega}$.
The existence and uniqueness of the solution are given by the following theorem:
Theorem 1. [10] Let $u_{0} \in L^{2}(\Omega)$ and $f \in L^{2}(] 0, T\left[; L^{2}(\Omega)\right)$, there is only one weak solution of the problem (1) such as:

$$
u \in L^{2}(] 0, T\left[; H_{0}^{1}(\Omega)\right) \cap C\left([0, T] ; L^{2}(\Omega)\right)
$$

with the following estimation:

$$
\begin{equation*}
\|u\|_{L^{2}(] 0, T\left[; H_{0}^{1}(\Omega)\right)}+\|u\|_{C\left([0, T] ; L^{2}(\Omega)\right)} \leq C\left(\|f\|_{L^{2}(] 0, T\left[; L^{2}(\Omega)\right)}+\left\|u_{0}\right\|_{L^{2}(\Omega)}\right) \tag{6}
\end{equation*}
$$

the constant $C$ depends only on $\Omega$ and $T$.
Theorem 2. Let $u$ a weak solution of (1) with $P=\left(a, u_{0}, f\right)$.
The function

$$
\begin{aligned}
\phi: S & \longrightarrow L^{2}(] 0, T\left[; H_{0}^{1}(\Omega)\right) \cap C\left([0, T] ; L^{2}(\Omega)\right) \\
& P \longmapsto u
\end{aligned}
$$

is continuous, and the functional $J_{T}$ is continuous. Therefore, $J_{T}$ has a unique minimum in $A_{a d} \square$
Proof. Let $\delta P=\left(\delta a, \delta u_{0}, \delta f\right) \in S$ be a small variation such that $P+\delta P \in A_{a d}$. Consider $\delta u=u^{\delta}-u$, where $u$ is the weak solution of (1) with $P$, and $u^{\delta}$ is the weak solution of (1) with $P^{\delta}=P+\delta P$. We have

$$
\left\{\begin{array}{l}
\left.\partial_{t} u-\partial_{x}\left(a \partial_{x} u\right)=f(x, t), \forall x \in \Omega, \forall t \in\right] 0, T[  \tag{7}\\
u(0, t)=u(1, t)=0, \forall t \in] 0, T[ \\
u(x, 0)=u_{0}(x), \quad \forall x \in \Omega
\end{array}\right.
$$

and

$$
\left\{\begin{array}{l}
\left.\partial_{t} u^{\delta}-\partial_{x}\left(a^{\delta} u_{x}^{\delta}\right)=f^{\delta}(x, t), \forall x \in \Omega, \forall t \in\right] 0, T[  \tag{8}\\
\left.u^{\delta}(0, t)=u^{\delta}(1, t)=0, \forall t \in\right] 0, T[ \\
u^{\delta}(x, 0)=u_{0}^{\delta}(x), \forall x \in \Omega
\end{array}\right.
$$

(8)-(7) yield

$$
\left\{\begin{array}{l}
\left.\partial_{t} \delta u-\partial_{x}\left(a^{\delta} \partial_{x} \delta u\right)-\partial_{x}\left(\delta a \partial_{x} u\right)=\delta f, \forall x \in \Omega, \forall t \in\right] 0, T[  \tag{9}\\
\delta u(0, t)=\delta u(1, t)=0, \forall t \in] 0, T[ \\
\delta u(x, 0)=\delta u_{0}(x), \forall x \in \Omega
\end{array}\right.
$$

The variational formulation of (9) is:

$$
\left\{\begin{array}{l}
\int_{\Omega} \partial_{t} \delta u v d x+\int_{\Omega}\left(a^{\delta} \partial_{x} \delta u\right) \partial_{x} v d x-\int_{\Omega} \partial_{x}\left(\delta a \partial_{x} u\right) v d x=\int_{\Omega} \delta f v d x, \forall v \in H_{0}^{1}(\Omega) \\
\delta u(x, 0)=\delta u_{0}(x), \forall x \in \Omega
\end{array}\right.
$$

Take $v=\delta u$, then

$$
\int_{\Omega} \partial_{t}(\delta u) \delta u d x+\int_{\Omega} a^{\delta}\left(\partial_{x} \delta u\right)^{2} d x-\int_{\Omega} \partial_{x}\left(\delta a \partial_{x} u\right) \delta u d x=\int_{\Omega} \delta f \delta u d x
$$

We have

$$
\int_{\Omega} a^{\delta}\left(\partial_{x} \delta u\right)^{2} d x \geq 0
$$

this implies that

$$
\int_{\Omega} \partial_{t}(\delta u) \delta u d x \leqslant \int_{\Omega} \partial_{x}\left(\delta a \partial_{x} u\right) \delta u d x+\int_{\Omega} \delta f \delta u d x
$$

and consequently

$$
\int_{\Omega} \partial_{t}(\delta u) \delta u d x \leqslant \int_{\Omega} \partial_{x}\left(\delta a \partial_{x} u\right) \delta u d x+\int_{\Omega} \delta f \delta u d x
$$

then

$$
\int_{\Omega} \partial_{t}(\delta u) \delta u d x \leqslant\|\delta a\|_{L^{\infty}(\Omega)} \int_{\Omega}\left|\partial_{x} u \partial_{x} \delta u\right| d x+\int_{\Omega}|\delta f \delta u| d x
$$

By integrating between 0 and t with $\mathrm{t} \in[0, T]$ we obtain

$$
\frac{1}{2}\|\delta u(t)\|_{L^{2}(\Omega)}^{2}-\frac{1}{2}\left\|\delta u_{0}\right\|_{L^{2}(\Omega)}^{2} \leqslant\|\delta a\|_{L^{\infty}(\Omega)} \int_{0}^{T} \int_{\Omega}\left|\partial_{x} u \partial_{x} \delta u\right| d x d t+\int_{0}^{T} \int_{\Omega}|\delta f \delta u| d x
$$

since $u, \delta u \in H^{1}\left(0, T ; L^{2}(\Omega)\right)$,
there exist $C_{1}>0$ and $C_{2}>0$, such that

$$
\sup _{t \in[0, T]}\|\delta u(t)\|_{L^{2}(\Omega)}^{2} \leqslant 2 C_{1}\|\delta a\|_{L^{\infty}(\Omega)}+\left\|\delta u_{0}\right\|_{L^{2}(\Omega)}^{2}+C_{2}\|\delta f\|_{L^{2}(] 0, T\left[; L^{2}(\Omega)\right)}
$$

which give

$$
\|\delta u\|_{C\left([0, T] ; L^{2}(\Omega)\right)}^{2} \leqslant C\left(\|\delta a\|_{L^{2}(\Omega)}+\left\|\delta u_{0}\right\|_{L^{2}(\Omega)}^{2}+\|\delta f\|_{L^{2}(] 0, T\left[; L^{2}(\Omega)\right)}\right)
$$

where $C=\max \left(2 C_{1}, 1, C_{2}\right)$.
We have:

$$
\begin{align*}
& H^{1}(\Omega) \underset{\text { compact }}{\hookrightarrow} L^{2}(\Omega),  \tag{10}\\
& H_{0}^{1}(\Omega) \underset{\text { compact }}{\hookrightarrow} L^{2}(\Omega) \tag{11}
\end{align*}
$$

Let

$$
F=\left\{g \in L^{2}\left([0, T] ; H^{1}(\Omega)\right) / \partial_{t} g \in L^{2}\left([0, T] ; L^{2}(\Omega)\right)\right\}
$$

By Aubin-Lions Lemma [11], :

$$
\begin{equation*}
F \underset{\text { compact }}{\longrightarrow} L^{2}\left([0, T] ; L^{2}(\Omega)\right) \tag{12}
\end{equation*}
$$

Let

$$
\begin{gathered}
A_{u_{0}}^{r_{1}}=\left\{u_{0} \in H_{0}^{1}(\Omega) ;\left\|u_{0}\right\|_{H_{0}^{1}(\Omega)} \leq r_{1}\right\}, \\
A_{a}^{r_{2}}=\left\{a \in H^{1}(\Omega) ;\|a\|_{H^{1}(\Omega)} \leq r_{2}\right\} \\
A_{f}^{r_{3}}=\left\{f \in F ;\|f\|_{L^{2}\left([0, T] ; H^{1}(\Omega)\right)} \leq r_{3}\right\},
\end{gathered}
$$

where $r_{1}, r_{2}, r_{3} \in \mathbb{R}^{*+}$. Let

$$
A_{a d}=A_{a}^{r} \times A_{u_{0}}^{r} \times A_{f}^{r},
$$

where $r=\min \left(r_{1}, r_{2}, r_{3}\right)$.
$A_{a}^{r}, A_{u_{0}}^{r}$ and $A_{f}^{r}$ are bounded in $H^{1}(\Omega), H_{0}^{1}(\Omega)$ and $F$ respectively. From, (11) and (12), they are compacts in $L^{2}(\Omega), L^{2}(\Omega)$ and $L^{2}\left([0, T] ; L^{2}(\Omega)\right)$ respectively. Consequently $A_{\text {ad }}$ is compact in $S$.

The functional $J_{T}$ is continuous. Therefore, $J$ has a unique minimum in $A_{\text {ad }}$.

## 4. Gradient of $J_{T}$ by adjoint state method

We define the Gâteaux derivative of $\phi$ at $P=\left(a, u_{0}, f\right)$ in the direction $q=\left(q_{a}, q_{u_{0}}, q_{f}\right) \in S$, by

$$
\hat{u}=\lim _{\varepsilon \rightarrow 0} \frac{\phi(P+\varepsilon q)-\phi(P)}{\varepsilon}
$$

$\phi(P+\varepsilon q)$ is the weak solution of (1) with parameters $P+\varepsilon q$, and $\phi(P)$ is the weak solution of (1) with parameters $P$.

The Gâteaux (directional) derivative of $\phi$ at $P$ in the direction $q \in S$, is the weak solution of the following tangent linear model:

$$
\left\{\begin{array}{l}
\left.\partial_{t} \hat{u}-\partial_{x}\left(a \partial_{x} \hat{u}\right)-\partial_{x}\left(q_{a} \partial_{x} u\right)=q_{f}(x, t), \forall x \in \Omega, \forall t \in\right] 0, T[, \\
\hat{u}(t, 0)=\hat{u}(t, 1)=0, \forall t \in] 0, T[ \\
\hat{u}(0, x)=q_{u_{0}}(x), \forall x \in \Omega .
\end{array}\right.
$$

We introduce the adjoint variable $A$, and we integrate:

$$
\int_{0}^{T} \int_{0}^{1} \partial_{t} \hat{u} A-\partial_{x}\left(a \partial_{x} \hat{u}\right) A-\partial_{x}\left(q_{a} \partial_{x} u\right) A d x d t=\int_{0}^{T} \int_{0}^{1} q_{f} A d x d t
$$

Calculate separately each term :

$$
\begin{gathered}
\int_{0}^{T} \int_{0}^{1} \partial_{t} \hat{u} A d x d t=\int_{0}^{1}[\hat{u} A]_{0}^{T} d x-\int_{0}^{1} \int_{0}^{T} \partial_{t} A \hat{u} d t d x \\
\begin{aligned}
\int_{0}^{T} \int_{0}^{1} \partial_{x}\left(a \partial_{x} \hat{u}\right) A d x d t & =\int_{0}^{T}\left[\left(a \partial_{x} \hat{u}\right) A\right]_{0}^{1} d t-\int_{0}^{T} \int_{0}^{1} a \partial_{x} \hat{u} \partial_{x} A d x d t \\
& =\int_{0}^{T}\left[\left(a \partial_{x} \hat{u}\right) A\right]_{0}^{1} d t-\int_{0}^{T}\left[a \hat{u} \partial_{x} A\right]_{0}^{1} d t+\int_{0}^{T} \int_{0}^{1} \partial_{x}\left(a \partial_{x} A\right) \hat{u} d x d t \\
& =\int_{0}^{T}\left[\left(a \partial_{x} \hat{u}\right) A\right]_{0}^{1} d t+\int_{0}^{T} \int_{0}^{1} \partial_{x}\left(a \partial_{x} A\right) \hat{u} d x d t
\end{aligned} \\
\int_{0}^{T} \int_{0}^{1} A\left(\partial_{x}\left(q_{a} \partial_{x} u\right)\right) d x d t=\int_{0}^{T}\left[q_{a} \partial_{x} u A\right]_{0}^{1} d t-\int_{0}^{T} \int_{0}^{1} q_{a} \partial_{x} u \partial_{x} A d x d t
\end{gathered}
$$

We pose $A(t, 0)=A(t, 1)=0$ and $A(T, x)=0$, we obtain

$$
\int_{0}^{T}\left\langle\hat{u}, \partial_{t} A+\partial_{x}\left(a \partial_{x} A\right)\right\rangle_{L^{2}(\Omega)} d t=\left\langle q_{a}, \int_{0}^{T} \partial_{x} u \partial_{x} A d t\right\rangle_{L^{2}(\Omega)}-\left\langle q_{u_{0}}, A(0, x)\right\rangle_{L^{2}(\Omega)}-\left\langle q_{f}, A\right\rangle_{L^{2}(00, T[\times \Omega)},
$$

so

$$
\begin{equation*}
\int_{0}^{T}\left\langle\partial_{t} A+\partial_{x}\left(a \partial_{x} A, \hat{u}\right)\right\rangle_{L^{2}(\Omega)} d t=-\langle q, h\rangle_{S} \tag{13}
\end{equation*}
$$

where

$$
h=\left(-\int_{0}^{T} \partial_{x} u \partial_{x} A d t, A(0, x), A\right)
$$

Let

$$
J_{T}(P)=\frac{1}{2 T} \int_{0}^{T}\left\|u(t)-u^{o b s}(t)\right\|_{L^{2}(\Omega)}^{2} d t
$$

The Gâteaux derivative of $J$ at $P$ in the direction $q$ is given by:

$$
\hat{J}_{P}(q)=\lim _{\varepsilon \rightarrow 0} \frac{J_{T}(P+\varepsilon q)-J_{T}(P)}{\varepsilon}
$$

After some computations, we arrive at

$$
\begin{equation*}
\hat{J}_{P}(q)=\frac{1}{T} \int_{0}^{T}\left\langle u-u^{o b s}, \hat{u}\right\rangle_{L^{2}(\Omega)} d t \tag{14}
\end{equation*}
$$

The adjoint model is:

$$
\left\{\begin{array}{l}
\partial_{t} A+\partial_{x}\left(a \partial_{x} A\right)=\frac{1}{T}\left(u-u^{o b s}\right)  \tag{15}\\
A(t, 0)=A(t, 1)=0, \forall t \in] 0, T[ \\
A(T, x)=0, \forall x \in \Omega
\end{array}\right.
$$

the problem is retrograde, we make the change of variable $t \longleftrightarrow T-t$, then

$$
\left\{\begin{array}{l}
\partial_{t} A-\partial_{x}\left(a \partial_{x} A\right)=\frac{1}{T}\left(\bar{u}^{o b s}-\bar{u}\right),  \tag{16}\\
A(t, 0)=A(t, 1)=0, \quad \forall t \in] 0, T[ \\
A(0, x)=0, \quad \forall x \in \Omega
\end{array}\right.
$$

with $\bar{u}(t)=u(T-t)$.
From equations (13), (14) et (15), the gradient of $J_{T}$ is given by

$$
\begin{equation*}
\nabla_{P} J_{T}=\left(\frac{\partial J_{T}}{\partial a}, \frac{\partial J_{T}}{\partial u_{0}}, \frac{\partial J_{T}}{\partial f}\right)=-h \tag{17}
\end{equation*}
$$

By the change of variable we have:

$$
\frac{\partial J_{T}}{\partial a}=\int_{0}^{T} \partial_{x} u \partial_{x} A d t ; \quad \frac{\partial J_{T}}{\partial u_{0}}=-A(T, x) ; \quad \frac{\partial J_{T}}{\partial f}=-A ;
$$

## 5. Approach, Model and Algorithm

Consider a parabolic PDE with $d$ spatial dimensions:

$$
\left\{\begin{array}{l}
\left.\partial_{t} u-\operatorname{div}_{x}\left(a(x) I_{d} \nabla_{x} u\right)=f, \forall x \in \Omega, \forall t \in\right] 0, T[  \tag{18}\\
\left.\left.u\right|_{\partial \Omega}=g(t), \forall t \in\right] 0, T[ \\
u(0, x)=u_{0}(x), \forall x \in \Omega
\end{array}\right.
$$

where $\Omega \subset \mathbb{R}^{d}, I_{d}$ is the identity matrix of order $d$.
The goal is to find $P=\left(a, u_{0}, f\right) \in S_{P}$ such that $U^{o b} \in \mathbb{R}^{k}$ are $k$ final observations of $u$ solution of (18) at $K$, with

$$
\begin{gathered}
S_{P}=L^{2}(\Omega) \times L^{2}(\Omega) \times L^{2}(] 0, T[\times \Omega), \\
K=\left\{\left(t_{i}, x_{i}\right) \in\{T\} \times \Omega\right\}_{0 \leqslant i \leqslant k}
\end{gathered}
$$

and $P^{b}=\left(a^{b}, u_{0}^{b}, f^{b}\right) \in \mathbb{R}^{h}$ is $h$ a priori knowledge of $P$ at $H^{b}$ with some bruit $b$ where $h=\left(h_{1}, h_{2}, h_{3}\right) \in \mathbb{N}^{3}$ :

$$
H^{b}=\left\{x_{i} \in \Omega\right\}_{0 \leqslant i \leqslant h_{1}} \times\left\{x_{i} \in \Omega\right\}_{0 \leqslant i \leqslant h_{2}} \times\left\{\left(t_{i}, x_{i}\right) \in\right] 0, T[\times \Omega\}_{0 \leqslant i \leqslant h_{3}}
$$

The optimization problem associated is finding $P$ which minimizes the function:

$$
\begin{equation*}
\mathcal{J}(v)=\left\|u(t=T)-u^{o b}\right\|_{L^{2}(\Omega)}^{2}+\left\|v-p^{b}\right\|_{S_{P}}^{2} \tag{19}
\end{equation*}
$$

where $u^{o b}$ and $p^{b}$ are interpolations of $U^{o b}$ in $L^{2}(\Omega)$ and $P^{b}$ in $S_{P}$.

The deep learning method for inverse problem (DL-IP) approach approximates $a, u_{0}, f$ and $u$ by four deep neural networks with common inputs and loss function. Figure (1) shows a visualization of the overall architecture:

$$
\begin{gathered}
a(x) \simeq \tilde{a}(x)=N_{a}\left(x ; \theta_{a}\right), \\
u_{0}(x) \simeq \tilde{u_{0}}(x)=N_{u_{0}}\left(x ; \theta_{u_{0}}\right), \\
f(x, t) \simeq \tilde{f}(x, t)=N_{f}\left(x, t ; \theta_{f}\right), \\
u(x, t) \simeq \tilde{u}(x, t)=N_{u}\left(x, t ; \theta_{u}\right),
\end{gathered}
$$

where $\theta=\left(\theta_{a}, \theta_{u_{0}}, \theta_{f}, \theta_{u}\right)$ are parameters of neural networks.


Figure 1. DL-IP's model architecture.

Our goal now is to find $\theta$ such that $U^{o b}$ are observations of $N_{u}$, and $P^{b}$ is the a priori knowledge of $\left(N_{a}, N_{u_{0}}, N_{f}\right)$.

The cost function is constructed as follows:

$$
\begin{align*}
& \tilde{\mathcal{J}}(\theta)=\left\|\partial_{t} \tilde{u}+\mathcal{L}(\tilde{a}, \tilde{u})-\tilde{f}\right\|_{L^{2}(] 0, T[\times \Omega)}^{2}+\left\|\tilde{u}^{o b}-U^{o b}\right\|_{\mathbb{R}^{k}}^{2}+\left\|\tilde{P}^{b}-P^{b}\right\|_{\mathbb{R}^{h}}^{2}+  \tag{20}\\
&\left\|\left.\tilde{u}\right|_{\partial \Omega}-g\right\|_{L^{2}(] 0, T[\times \partial \Omega)}^{2}+\left\|\tilde{u}(t=0)-\tilde{u_{0}}\right\|_{L^{2}(\Omega)}^{2},
\end{align*}
$$

where $\mathcal{L}(a, u)=-\operatorname{div}_{x}\left(a(x) I_{d} \nabla_{x} u\right), \tilde{u}^{o b}$ are values of $\tilde{u}$ at $K, \tilde{P^{b}}$ are values of $\left(\tilde{a}, \tilde{u_{0}}, \tilde{f}\right)$ at $H^{b}$.
The loss functional $\tilde{\mathcal{J}}$ consists of five parts:

1. A measure of how well $\tilde{\mathcal{U}}$ and $\tilde{P}$ satisfy the differential operator:

$$
\tilde{\mathcal{J}}_{1}(\theta)=\left\|\partial_{t} \tilde{u}+\mathcal{L}(\tilde{a}, \tilde{u})-\tilde{f}\right\|_{L^{2}(] 0, T[\times \Omega)}^{2}
$$

2. A measure of how well the approximation $\tilde{u}$ satisfies the observation value:

$$
\tilde{\mathcal{J}}_{2}(\theta)=\left\|\tilde{u}^{o b}-U^{o b}\right\|_{\mathbb{R}^{k}}^{2}
$$

3. A measure of how well $\tilde{P}$ satisfies the a priori knowledge:

$$
\tilde{\mathcal{J}}_{3}(\theta)=\left\|\tilde{P}-P^{b}\right\|_{\mathbb{R}^{h}}^{2}
$$

4. A measure of how well the approximation $\tilde{u}$ satisfies the boundary condition:

$$
\tilde{\mathcal{J}}_{4}(\theta)=\left\|\left.\tilde{u}\right|_{\partial \Omega}-g\right\|_{L^{2}(] 0, T[\times \partial \Omega)}^{2}
$$

5. A measure of how well the approximation $\tilde{u}$ satisfies the approximation of initial condition $\tilde{u_{0}}$ :

$$
\tilde{\mathcal{J}}_{5}(\theta)=\left\|\tilde{u}(t=0)-\tilde{u_{0}}\right\|_{L^{2}(\Omega)}^{2}
$$

The derivatives of $\tilde{u}$ can be evaluated using automatic differentiation (see [12]) since it is parameterizing as a neural network.

The architecture of a neural network can be crucial to its success. Frequently, different applications require different architectures. For example, convolution networks are essential for image recognition while long short-term networks (LSTMs) are useful for modeling sequential data. In DL-IP's model presented by Figure (1), we use four neural networks with the same architecture described in figure (2).


Figure 2. The structure of networks used in figure (1).

Each neural network $N_{k \in\left\{a, u_{0}, f, u\right\}}$ uses three types of layers $L_{1}, L_{2}$ and $L_{3}$.

- $L_{1}$ takes as input $I=(x, t)$ or $I=x$ and returns as output $C^{0}$ :

$$
C^{0}=\sigma\left(\theta^{0} . I+\beta^{0}\right)
$$

- $L_{2}^{n-1}$ takes as inputs $I=(x, t)$ or $I=x$ and the output $C^{n-1}$ of the previous layer $L_{1}$ if $n=1$ or $L_{2}^{n-2}$ if $n>1$ and returns as output $C^{n}$. The inputs are transformed through a series of operations. Below, we present the architecture in the equations along with a visual representation of a single $L_{2}^{n-1}$ layer in Figure (3):

$$
\begin{align*}
Z_{1}^{n-1} & =\sigma\left(\theta_{1}^{n} \cdot I+\gamma_{1}^{n} \cdot C^{n-1}+\beta_{1}^{n}\right) \\
Z_{2}^{n-1} & =\sigma\left(\theta_{2}^{n} \cdot I+\gamma_{2}^{n} \cdot C^{n-1}+\beta_{2}^{n}\right) \\
Z_{3}^{n-1} & =\sigma\left(\theta_{3}^{n} \cdot I+\gamma_{3}^{n} \cdot C^{n-1}+\beta_{3}^{n}\right)  \tag{21}\\
Z_{4}^{n-1} & \left.=\theta_{4}^{n} \cdot I+\gamma_{4}^{n} \cdot\left(C^{n-1} \odot Z_{1}^{n-1}\right)+\beta_{4}^{n}\right), \\
C^{n} & =\left(1-Z_{2}^{n-1}\right) \odot Z_{4}^{n-1}+Z_{3}^{n-1} \odot C^{n-1},
\end{align*}
$$

where $1 \leqslant n \leqslant L$, $\odot$ denotes Hadamard multiplication, $\theta^{n}=\left\{\theta_{i}^{n}, \gamma_{i}^{n}, \beta_{i}^{n}\right\}_{1 \leqslant i \leqslant 4}$ are parameters of $L_{2}^{n-1}$ layer and $L$ is the total number of $L_{2}$ layers.


Figure 3. Operations within layer $L_{2}^{n-1}$.

- $L_{3}$ takes as input the output $C^{L}$ of the last layer $L_{2}^{L}$ and returns:

$$
C^{L+1}=\theta^{L+1} . I+\beta^{L+1}
$$

$\sigma$ is an activation function such as the tanh function or sigmoidal function, $\theta_{k}$ are parameters of the $N_{k}$ neural network and $\theta$ are parameters of the DL-IP's model, where

$$
\begin{gathered}
\theta_{k}=\left\{\theta^{0}, \beta^{0}, \ldots, \theta^{n}, \ldots, \theta^{L+1}, \beta^{L+1}\right\}_{1 \leqslant n \leqslant L} \\
\theta=\left\{\theta_{k}\right\}_{k \in\left\{a, u_{0}, f, u\right\}}
\end{gathered}
$$

## Algorithm 3. DL-IP Algorithm

## 1. Define boundary conditions

2. Define the architecture of neutral networks by setting the number of layers, number of neurons in each layer and activation functions.
3. Generate random training set:

$$
D_{M}=\left\{\left(t_{i}, x_{i}\right)_{0 \leqslant i \leqslant \tilde{M}} \in[0, T] \times \Omega,\left(\overline{t_{i}}, \overline{x_{i}}\right)_{0 \leqslant i \leqslant \bar{M}} \in[0, T] \times \partial \Omega,\left(x_{i}^{0}\right)_{0 \leqslant i \leqslant M^{0}} \in \Omega\right\},
$$

where $M=\left(\tilde{M}, \bar{M}, M^{0}\right) \in \mathbb{N}^{* 3}$.
4. Initialize the parameter set $\theta_{0}$ and the learning rate $\alpha_{0}$.
5. Repeat until convergence criterion is satisfied.

1 Randomly sample a mini-batch $d_{m}$ of training examples from $D_{M}$ :

$$
d_{m}=\left\{\left(t_{i}, x_{i}\right)_{0 \leqslant i \leqslant \tilde{m} \leqslant \tilde{M}},\left(\overline{t_{i}}, \overline{x_{i}}\right)_{0 \leqslant i \leqslant \bar{m} \leqslant \bar{M}},\left(x_{i}^{0}\right)_{0 \leqslant i \leqslant m^{0} \leqslant M^{0}}\right\} \subset D_{M} .
$$

2 Compute the loss functional for the sampled mini-batch $d_{m}$ :

$$
\begin{gather*}
\tilde{J}_{1}^{m}\left(\theta_{n}\right)=\frac{1}{m} \sum_{0}^{m}\left(\partial_{t} \tilde{u}\left(t_{i}, x_{i}\right)+\mathcal{L}\left(\tilde{a}\left(t_{i}, x_{i}\right), \tilde{u}\left(t_{i}, x_{i}\right)\right)-\tilde{f}\left(t_{i}, x_{i}\right)\right)^{2}, \\
\tilde{J}_{2}^{k}\left(\theta_{n}\right)=\frac{1}{k} \sum_{\left(t_{i}, x_{i}\right) \in K}\left(\tilde{u}\left(t_{i}, x_{i}\right)-U_{i}^{o b}\right)^{2}, \\
\tilde{J}_{3}^{h}\left(\theta_{n}\right)=\frac{1}{\bar{h}} \sum_{\left(t_{i}, x_{i}\right) \in B}\left(\left(\tilde{a}\left(x_{i}\right)-a_{i}^{b}\right)^{2}+\left(\tilde{f}\left(t_{i}, x_{i}\right)-f_{i}^{b}\right)^{2}+\left(\tilde{u_{0}}\left(t_{i}, x_{i}\right)-u_{0 i}^{b}\right)^{2}\right), \\
\tilde{J}_{4}^{\bar{m}}\left(\theta_{n}\right)=\frac{1}{\bar{m}} \sum_{i=0}^{\bar{m}}\left(\tilde{u}\left(\overline{t_{i}}, \overline{x_{i}}\right)-g\left(\overline{x_{i}}\right)\right)^{2}, \\
\tilde{J}_{5}^{m^{0}}\left(\theta_{n}\right)=\frac{1}{m^{0}} \sum_{i=0}^{m^{0}}\left(\tilde{u}\left(0, x_{i}^{0}\right)-\tilde{u_{0}}\left(x_{i}^{0}\right)\right)^{2}, \\
\tilde{J}\left(\theta_{n}, d_{m}\right)=\tilde{J}_{1}^{m}\left(\theta_{n}\right)+\tilde{J}_{2}^{k}\left(\theta_{n}\right)+\tilde{J}_{3}^{h}\left(\theta_{n}\right)+\tilde{J}_{4}^{m^{0}}\left(\theta_{n}\right)+\tilde{J}_{5}^{m}\left(\theta_{n}\right) . \tag{22}
\end{gather*}
$$

3 Compute the gradient $\nabla_{\theta_{n}} \tilde{J}\left(\theta_{n}, d_{m}\right)$ for the sampled mini-batch $d_{m}$ using backpropagation.
4 Use the estimated gradient to take a descent step at $d_{m}$ with learning rates (23) to update $\theta_{n+1}$ :

$$
\begin{gather*}
\theta_{n+1}=\theta_{n}-\alpha_{n} \nabla_{\theta_{n}} \tilde{J}\left(\theta_{n}, d_{m}\right) \\
\alpha_{n}=\left\{\begin{array}{l}
10^{-4}, \quad n \leqslant 5000 \\
5 \times 10^{-4}, \quad 5<n \leqslant 10000 \\
10^{-5}, \quad 10<n \leqslant 20000 \\
5 \times 10^{-5}, \quad 20<n \leqslant 30000, \\
10^{-6}, \quad 5<30 \leqslant 40000 \\
5 \times 10^{-6}, \quad 40<n \leqslant 50000 \\
10^{-7}, \quad 50000<n
\end{array}\right. \tag{23}
\end{gather*}
$$

6. Save model to be used for any $x \in \Omega$ and $t \in] 0, T[$.

## 6. Convergence of the $L^{2} \operatorname{error} \tilde{\mathcal{J}}(\tilde{u}, \tilde{P})$

Let the $L^{2} \operatorname{error} \tilde{\mathcal{J}}(\tilde{u}, \tilde{P})$ measure how well neural networks $\tilde{u}$ and $\tilde{P}=\left(\tilde{a}, \tilde{u_{0}}, \tilde{f}\right)$ satisfy observations, differential operator, boundary condition, initial condition and a priori knowledge.

$$
\begin{align*}
\tilde{\mathcal{J}}(\tilde{u}, \tilde{P})= & \left\|\partial_{t} \tilde{u}+\mathcal{L}(\tilde{a}, \tilde{u})-\tilde{f}\right\|_{L^{2}(] 0, T[\times \Omega)}^{2}+\left\|\tilde{u}(t=T)-u^{o b}\right\|_{L^{2}(\Omega)}^{2}+\left\|\tilde{P}-p^{b}\right\|_{S_{P}}^{2}+  \tag{24}\\
& \left\|\left.\tilde{u}\right|_{\partial \Omega}-g\right\|_{L^{2}(] 0, T[\times \partial \Omega)}^{2}+\left\|\tilde{u}(t=0)-\tilde{u_{0}}\right\|_{L^{2}(\Omega)}^{2}
\end{align*}
$$

where $u^{o b}$ and $p^{b}$ being interpolations of $U^{o b}$ in $L^{2}(\Omega)$ and $P^{b}$ in $S_{P}$.
$\tilde{\mathcal{J}}(\tilde{u}, \tilde{P})$ can be directly calculated from the PDE (18) for any approximations $(\tilde{u}, \tilde{P})$. The goal is to construct functions $(\tilde{u}, \tilde{P})$ for which $\tilde{\mathcal{J}}(\tilde{u}, \tilde{P})$ is as close to 0 as possible. Define $\mathfrak{C}^{n}$ and $\mathfrak{C}_{\mathrm{t}}{ }^{n}$ as classes of neural networks with a single hidden layer and $n$ hidden units. Let $\left(\tilde{u}^{n}, \tilde{P}^{n}\right)$ be neural networks with $n$ hidden units which minimize $\tilde{\mathcal{J}}(\tilde{u}, \tilde{P})$. We prove that, under certain conditions,

$$
\text { there exists }\left(\tilde{u}^{n}, \tilde{P}^{n}\right) \in \mathfrak{C}_{4}^{n} \text { such that } \tilde{\mathcal{J}}\left(\tilde{u}^{n}, \tilde{P}^{n}\right) \rightarrow 0 \text {, as } n \rightarrow \infty .
$$

strongly in $L^{2}([0, T] \times \Omega) \times S$.
The proof requires the joint analysis of the approximation power of neural networks as well as the continuity properties of partial differential equations. We show that the neural networks can satisfy observations, the differential operator, boundary condition, and initial condition arbitrarily well for sufficiently large $n$,

$$
\tilde{\mathcal{J}}\left(\tilde{u}^{n}, \tilde{P}^{n}\right) \rightarrow 0 \text { as } n \rightarrow \infty
$$

In this section, we present a theorem guaranteeing the existence of multilayer feed forward networks $\left(\tilde{u}^{n}, \tilde{P}^{n}\right)$ able to universally approximate solutions $(u, P)$ in the sense that there is ( $\left.\tilde{u}^{n}, \tilde{P}^{n}\right)$ that makes the objective function $\tilde{\mathcal{J}}$ arbitrarily small. To do so, we use the results of [13] on universal approximation of functions and their derivatives. We denote

$$
\Omega_{T}=[0, T] \times \Omega \quad \mathcal{O}[u, P]=\partial_{t} u-\operatorname{div}_{x}\left(a(x) I_{d} \nabla_{x} u\right)-f .
$$

Universal approximation results for single functions and their derivatives have been obtained under various assumptions in [13-15]. In this paper, we use Theorem 3 of [13]. Let us recall the setup appropriately modified for our case of interest. Let $\psi$ be an activation function, e.g., of sigmoid type, of the hidden units and define the sets

$$
\begin{equation*}
\mathfrak{C}_{\mathrm{t}}^{n}(\psi)=\left\{\zeta_{t}(t, x): \mathbb{R}^{1+d} \mapsto \mathbb{R}: \zeta_{t}(t, x)=\sum_{i=1}^{n} \beta_{i}^{t} \psi\left(\alpha_{0, i}^{t} t+\sum_{j=1}^{d} \alpha_{j, i}^{t} x_{j}+c_{j}^{t}\right)\right\} \tag{25}
\end{equation*}
$$

where $\theta_{t}=\left(\beta_{1}^{t}, \ldots, \beta_{n}^{t}, \alpha_{0,1}^{t}, \ldots, \alpha_{d, n}^{t}, c_{1}^{t}, \ldots, c_{d}^{t}\right) \in \mathbb{R}^{2 n+d(1+n)}$ compose the elements of the parameter space, and

$$
\begin{equation*}
\mathfrak{C}^{n}(\psi)=\left\{\zeta(x): \mathbb{R}^{d} \mapsto \mathbb{R}: \zeta(x)=\sum_{i=1}^{n} \beta_{i} \psi\left(\sum_{j=1}^{d} \alpha_{j, i} x_{j}+c_{j}\right)\right\} \tag{26}
\end{equation*}
$$

where $\theta=\left(\beta_{1}, \ldots, \beta_{n}, \alpha_{1,1}, \ldots, \alpha_{d, n}, c_{1}, \ldots, c_{n}\right) \in \mathbb{R}^{n+d(1+n)}$ compose the elements of the parameter space. We consider the product spaces:

$$
\mathfrak{C}_{4}^{n}(\psi)=\mathfrak{C}_{\mathfrak{t}}^{n}(\psi) \times \mathfrak{C}^{n}(\psi) \times \mathfrak{C}^{n}(\psi) \times \mathfrak{C}_{\mathfrak{t}}^{n}(\psi) ; \quad \mathfrak{C}_{4}(\psi)=\mathfrak{C}_{\mathfrak{t}}(\psi) \times \mathfrak{C}(\psi) \times \mathfrak{C}(\psi) \times \mathfrak{C}_{\mathfrak{t}}(\psi)
$$

where

$$
\mathfrak{C}_{\mathfrak{t}}(\psi)=\bigcup_{n \geq 1} \mathfrak{C}_{\mathfrak{t}}^{n}(\psi) ; \quad \mathfrak{C}(\psi)=\bigcup_{n \geq 1} \mathfrak{C}^{n}(\psi)
$$

Then we have the following result.
Theorem 4. Let $\mathfrak{C}^{n}(\psi)$ and $\mathfrak{C}_{\mathrm{t}}{ }^{n}(\psi)$ be given by (26) and 25 respectively where $\psi$ is assumed to be in $C^{2}(\mathbb{R})$, bounded and non-constant. Assume that $\Omega_{T}$ is compact. In addition, assume that the problem has a unique solution $(u, P)$. Then, for every $\epsilon>0$, there exists a positive constant $K>0$ such that there exists a $(\tilde{u}, \tilde{P}) \in \mathfrak{C}_{4}(\psi)$ that satisfies

$$
\tilde{\mathcal{J}}(\tilde{u}, \tilde{P}) \leq K \epsilon
$$

## Proof.

By Theorem 3 of [13] we know that there is a function $(\tilde{u}, \tilde{P}) \in \mathfrak{C}_{4}(\psi)$ that is uniformly dense on compacts of $C_{4}$. Where $C_{4}=\mathcal{C}^{2}\left(\mathbb{R}^{1+d}\right) \times \mathcal{C}^{2}\left(\mathbb{R}^{d}\right) \times \mathcal{C}\left(\mathbb{R}^{d}\right) \times \mathcal{C}\left(\mathbb{R}^{1+d}\right)$. This means that for $\left(u, a, u_{0}, f\right) \in C_{4}$ and $\epsilon>0$, there is $\left(\tilde{u}, \tilde{a}, \tilde{u_{0}}, \tilde{f}\right) \in \mathfrak{C}_{4}(\psi)$ such that

$$
\sup _{(t, x) \in \Omega_{T}}\left|\partial_{t} u(t, x)-\partial_{t} \tilde{u}\left(t, x ; \theta_{u}\right)\right|+\max _{|\alpha| \leq 2} \sup _{(t, x) \in \Omega_{T}}\left|\partial_{x}^{(\alpha)} u(t, x)-\partial_{x}^{(\alpha)} \tilde{u}\left(t, x ; \theta_{u}\right)\right|<\epsilon
$$

$$
\begin{aligned}
& \max _{|\alpha| \leq 2} \sup _{x \in \bar{\Omega}}\left|\partial_{x}^{(\alpha)} a(x)-\partial_{x}^{(\alpha)} \tilde{a}\left(x ; \theta_{a}\right)\right|<\epsilon \\
& \sup _{x \in \bar{\Omega}}\left|u_{0}(x)-\tilde{u_{0}}\left(x ; \theta_{u_{0}}\right)\right|<\epsilon \\
& \sup _{(t, x) \in \Omega_{T}}\left|f(t, x)-\tilde{f}\left(t, x ; \theta_{f}\right)\right|<\epsilon
\end{aligned}
$$

For the objective function (note that $\mathcal{O}[u, P](t, x)=0$ for $(u, P)$ that solves the PDE)

$$
\begin{aligned}
\tilde{\mathcal{J}}(\tilde{u}, \tilde{P})= & \|\mathcal{O}[\tilde{u}, \tilde{P}]-\mathcal{O}[u, P]\|_{L^{2}\left(\Omega_{T}\right)}^{2}+\left\|\tilde{u}(t=T)-u^{o b}\right\|_{L^{2}(\Omega)}^{2}+\left\|\tilde{P}-p^{b}\right\|_{S_{P}}^{2}+ \\
& \left\|\left.\tilde{u}\right|_{\partial \Omega}-g\right\|_{L^{2}(] 0, T[\times \partial \Omega)}^{2}+\left\|\tilde{u}(t=0)-\tilde{u_{0}}\right\|_{L^{2}(\Omega)}^{2} \\
& \leq\left\|\partial_{t} u(t, x)-\partial_{t} \tilde{u}\left(t, x ; \theta_{u}\right)\right\|_{L^{2}\left(\Omega_{T}\right)}^{2}+\left\|f(t, x)-\tilde{f}\left(t, x ; \theta_{f}\right)\right\|_{L^{2}\left(\Omega_{T}\right)}^{2}+\|\mathcal{L}(a, u)-\mathcal{L}(\tilde{a}, \tilde{u})\|_{L^{2}\left(\Omega_{T}\right)}^{2}+ \\
& \left\|\tilde{u}(t=T)-u^{o b}\right\|_{L^{2}(\Omega)}^{2}+\left\|\tilde{f}\left(t, x ; \theta_{f}\right)-f^{b}(t, x)\right\|_{L^{2}\left(\Omega_{T}\right)}^{2}+\left\|\tilde{a}\left(x ; \theta_{a}\right)-a^{b}(x)\right\|_{L^{2}(\Omega)}^{2}+ \\
& \left\|\tilde{u_{0}}\left(x ; \theta_{u_{0}}\right)-u_{0}^{b}(x)\right\|_{L^{2}(\Omega)}^{2}+\left\|\left.\tilde{u}\right|_{\partial \Omega}-g\right\|_{L^{2}(] 0, T[\times \partial \Omega)}^{2}+\left\|\tilde{u}(t=0)-\tilde{u_{0}}\right\|_{L^{2}(\Omega)}^{2} \\
& \leq K \epsilon
\end{aligned}
$$

## 7. Numerical results and Analysis

### 7.1. One dimensional Heat equation

Let us start with the heat equation with Dirichlet boundary conditions. In one space dimension, the problem reads as:

$$
\left\{\begin{array}{l}
\left.\partial_{t} u-\partial_{x}\left(a \partial_{x} u\right)=f, \forall x \in\right] 0,1[, \forall t \in] 0,1[  \tag{27}\\
u(t, 0)=u(t, 1)=0, \forall t \in] 0,1[ \\
\left.u(0, x)=u_{0}(x), \forall x \in\right] 0,1[
\end{array}\right.
$$

From measures at 50 couples of $(T, x)$ created from 50 space points at final instant, we try to find $P=\left(a, u_{0}, f\right)$. We represent the solution $u$ and $P$ by four neural networks with $5 L_{2}$ layers and 10 neurons in each sub-layer. As for the activation functions, we use tanh. In general, the choice of a neural network's architecture (number of layers, neurons and form of activation functions) is crucial and in many cases still remains an art that relies on one's ability to balance the trade off between expressivity and trainability of the neural network. Our tests indicate that deeper and wider networks are usually more expressive but are often more costly to train (i.e., a feed-forward evaluation of the neural network takes more time and the optimizer requires more iterations to converge). The neural networks are trained by minimizing the sum of squared errors loss of equation (22).

To illustrate the effectiveness of our approach, we applied method of manufactured solution (MMS), i.e. construct some equations with analytic solutions and then compare the numerical solution with the analytical solutions..

Consider the analytical solutions of (27):

$$
a(x)=\frac{1+e^{x}}{4} ; \quad u(t, x)=\frac{\left(e^{x}-1\right)(x-1)^{2}}{1+t}
$$

from $a$ and $u$ below, we can calculate the analytical source $f$ and initial condition $u_{0}$. The neural network parameters are initialized using normal initialization and updated using the well-known ADAM algorithm (see [16]) with a decaying learning rate schedule (23). We use approximately 100 epochs. The data set of training is generated randomly by uniform distribution. 90 time points and 90 space points are chosen
randomly to create 8100 couples of $(t, x)$ to train the neural networks. The data is sampled to mini-batches of size 96. We implement the algorithm using TensorFlow which is a software library for deep learning.


Figure 4. Data generation for training neural networks. Here • represent the sampling points within the domain, - represent the sampling points on the boundary, - represent the initial sampling points and $\times$ represent the points of observations.

After training, we test the trained model on a testing data of 2916 couples of $(t, x)$ which are uniformly generated. The results below are calculated on the testing data.


Figure 5. The exact solution $u$ of equation (27) (left panel) is compared to the solution $\tilde{u}$ predicted by DL-IP (right panel). It is hard to distinguish predicted solution and exact solution by eyes.


Figure 6. Absolute space errors between predicted solution and exact solution at all points of time. The maximal error value is $8.3997 e-04$. On the training data points it is $2.0591 e-06$.


Figure 7. The exact coefficient $a$ is presented by - and the coefficient $\tilde{a}$ predicted by DL-IP is presented by - (left panel). Right panel represents the absolute error between the solutions. It is also hard to distinguish predicted and exact coefficients by eyes. The maximal error value is $2.0306 e-03$. On the training data points it is $4.3610 e-07$.


Figure 8. The exact initial condition $u_{0}$ is presented by - and the predicted one $\tilde{u_{0}}$ is presented by $\circ$ (left panel). Right panel represents the absolute error between the solutions. It is also hard to distinguish predicted and exact initial conditions by eyes. The maximal error value is $8.3997 e-04$. On the training data points it is $8.1583 e-08$.


Figure 9. The exact source $f$ of equation (27) (left panel) is compared to $\tilde{f}$ predicted by DL-IP (right panel). It is hard to distinguish exact and predicted sources by eyes.


Figure 10. Absolute space errors between exact and predicted sources at all time points of testing data. The maximal error value is $2.2345 e-02$. On the training data points it is $3.2232 e-05$.


Figure 11. Variation of functional $\tilde{J}$ within 100 epochs which are equivalent to 9200 steps. Across the iterations the loss is decreasing even if there are some oscillations; its minimal value is $3.1398 e-05$.

Furthermore, we performed two systematic studies of the impact of noise levels and a priori knowledge amounts by keeping the total number of training observations as well as the neural network architectures fixed to the settings described above. The results of the first study are summarized in table (1) and figure (12). For the second one we present the results in table( 2 ) and figure (13).

Table 1. Variation of functional $\tilde{J}$ in terms of $b$ the noise between real value of $P$ and the used one $P^{b}$, it means $P^{b}=P+b .(b=0)$ is corresponding to the clear data. Here, the a priori knowledge is kept fixed at 100 points.

| Noise $b$ | 0 | $10 e-5$ | $10 e-4$ | $10 e-2$ |
| :--- | :--- | :--- | :--- | :--- |
| $\tilde{J}$ | $4.9279 e-05$ | $3.8593 e-05$ | $4.6976 e-05$ | $5.8740 e-05$ |
| error $_{a}$ | $1.3595 e-03$ | $3.0966 e-03$ | $2.544641 e-03$ | $1.0210 e-01$ |
| error $_{u_{0}}$ | $6.5393 e-04$ | $6.3966 e-04$ | $1.419791 e-03$ | $1.0136 e-01$ |
| error $_{f}$ | $3.9217 e-02$ | $3.7173 e-02$ | $3.065276 e-02$ | $1.1334 e-01$ |
| error $_{u}$ | $6.5393 e-04$ | $2.8137 e-03$ | $1.4197 e-03$ | $1.0136 e-01$ |



Figure 12. Graph of table (1) results. By increasing noise $b$ minimal value of $\tilde{J}$ increases too.

Table 2. Variation of functional $\tilde{J}$ in terms of $h$ the amount of a priori knowledge. Here, the noise is kept fixed at 0.

| $h$ | 25 | 900 | 2500 | 3600 | 6400 |
| :--- | :--- | :--- | :--- | :--- | :--- |
| $\tilde{J}$ | $5.4914 e-05$ | $5.4113 e-05$ | $3.6522 e-05$ | $3.5627 e-05$ | $3.2917 e-05$ |
| error $_{a}$ | $6.9173 e-03$ | $1.8087 e-03$ | $1.6346 e-03$ | $1.0738 e-03$ | $3.4335 e-03$ |
| error $_{u_{0}}$ | $6.9173 e-03$ | $1.6645 e-03$ | $2.0992 e-04$ | $3.5799 e-04$ | $3.3736 e-04$ |
| error $_{f}$ | $1.5644 e-01$ | $5.8674 e-02$ | $6.7514 e-02$ | $5.3573 e-02$ | $5.1853 e-02$ |
| error $_{u}$ | $6.9173 e-03$ | $1.1827 e-03$ | $4.7244 e-04$ | $2.5719 e-04$ | $5.7335 e-04$ |



Figure 13. Graph of table (2) results. By increasing the amount of a priori knowledge $h$ minimal value of $\tilde{J}$ decreases.

The key observation here is that less noise in the data enhances the performance of the algorithm. Our experience so far indicates that the negative consequences of more noise in the data can be remedied to some extent by obtaining more data.

### 7.2. Limitation of DL-IP method

For inverse or identification problem associated to the evolution problem like heat equation, it is interesting to reconstruct solutions from final observations and without a priori knowledge $h=0$. DL-IP method can reconstruct the solutions from final observations but with at least 25 of a priori knowledge. This is the main drawback associated with DL-IP method. We present in figures below some results of this limitation.


Figure 14. The solution $\tilde{u}$ predicted by DL-IP without any a priori knowledge. The maximal error value is 0.0423 .


Figure 15. The solution $\tilde{f}$ predicted by DL-IP without any a priori knowledge. The maximal error value is $3.4016 e+01$.


Figure 16. The exact coefficient $a$ is presented by - and the coefficient $\tilde{a}$ predicted by DL-IP without any a priori knowledge is presented by o . The maximal error value is $9.6417 e-01$


Figure 17. The exact initial condition $u_{0}$ is presented by - and the predicted one without any a priori knowledge $\tilde{u}_{0}$ is presented by o. The maximal error value is $2.9096 e-01$.

It should also be noted that from an a priori final knowledge the reconstruction cannot be done. Even if the coefficient $a$ and the initial condition $u_{0}$ are well predicted, the source term $f$ is poorly predicted since it is a spatio-temporal function.


Figure 18. The solution $\tilde{u}$ predicted by DL-IP with final a priori knowledge. The maximal error value is $7.2011 e-$ 02.


Figure 19. The solution $\tilde{f}$ predicted by DL-IP without any a priori knowledge. The maximal error value is $3.1105 e+00$.


Figure 20. The exact coefficient $a$ is presented by - and the coefficient $\tilde{a}$ predicted by DL-IP with final a priori knowledge is presented by $\circ$. The maximal error value is $1.1479 e-03$.


Figure 21. The exact initial condition $u_{0}$ is presented by - and the predicted one with final a priori knowledge $\tilde{u}_{0}$ is presented by o . The maximal error value is $3.0998 e-04$.

## 8. Comparison of methods

In this section we make a comparison between DL-IP's model and classical methods such as Tikhonov regularization and hybrid approach.

Firstly we make a general comparison of algorithms in table (3) by answering the following questions:

- Is the method meshfree (MF)?
- Does it require interpolations (I)?
- Does it need Numeric Integration (NI)?
- How many direct problems we need to solve (DP)?

Table 3. General comparison of algorithms.

|  | DL-IP model | Tikhonov regularization | Hybrid Method |
| :--- | :--- | :--- | :--- |
| MF | No | Yes | Yes |
| I | No | Yes | Yes |
| NI | No | Yes | Yes |
| DP | 0 | 2 by descent iteration | thead2 by descent iteration +1 by fitness call |

Secondly, we consider the example (7.1), we solve the problem by Tikhonov regularization (table 4), Hybrid Method (table 5), and we make a comparison with results of DL-IP's model (table 6) by considering amount of knowledge $h$, noise $b$, minimal values of $\tilde{J}$ or $J_{T}$ and errors between exact solution and predicted one (absolute error error $_{u}$ and relative error \%error ${ }_{u}$ ). We consider 50 observations, the size of mesh is $S=M N$ with $M=N=100$. Parameters of Armijo-Goldstein's rule are $\alpha=0.3, \beta=10 e-4, \alpha_{i}=0.9$. Regularization coefficient is $\varepsilon=0.005$. we turn hybrid approach with an initial population of size 10000,500 individuals and 500 generations.

## - Tikhonov regularization.

Table 4. Results of the Tikhonov approach. By increasing noise $b$ or decreasing a priori knowledge $h$, minimal values of $J_{T}$ and errors increase. Solution cannot be reconstructed from partial knowledge

| $h$ | $b$ | $J_{T}$ | error $_{u}$ | orrror $_{u}$ |
| :--- | :--- | :--- | :--- | :--- |
| 100 | $10 e-2$ | $2.3471 e-05$ | $3.1232 e-01$ | $1.8347 e+00$ |
|  | $10 e-5$ | $2.9642 e-05$ | $3.1758 e-01$ | $1.7483 e+00$ |
| 900 | $10 e-2$ | $1.8643 e-05$ | $3.0199 e-01$ | $2.0137 e+00$ |
|  | $10 e-5$ | $2.7292 e-05$ | $2.9533 e-01$ | $1.8559 e+00$ |
| $S$ | $10 e-2$ | $1.1090 e-06$ | $1.0162 e-01$ | $2.2864 e-01$ |
|  | $10 e-5$ | $1.0894 e-06$ | $3.6399 e-03$ | $2.4319 e-02$ |

## - Hybrid method

Table 5. Results given by the hybrid method with $b=10 e-4$. By decreasing a priori knowledge $h$, minimal values of $J_{T}$ and errors increase. Solution cannot be reconstructed from partial knowledge.

| $h$ | $J_{T}$ | error $_{u}$ | \%error $_{u}$ |
| :--- | :--- | :--- | :--- |
| 100 | $1.8332 e-05$ | $3.0238 e-01$ | $1.8283 e+00$ |
| 900 | $1.6953 e-05$ | $3.0160 e-01$ | $1.9380 e+00$ |
| $S$ | $1.9603 e-09$ | $2.3773 e-03$ | $2.2907 e-03$ |

## - DL-IP's model

Table 6. Results given by DL-IP's model.

| $h$ | $b$ | $\tilde{J}$ | error $_{u}$ | \%error $_{u}$ |
| :---: | :---: | :---: | :---: | :---: |
| 25 | $10 e-2$ | $2.7900 e-05$ | $1.1098 e-01$ | $2.3676 e-01$ |
|  | $10 e-5$ | $3.0141 e-05$ | $2.2934 e-02$ | $1.6950 e-02$ |
| 100 | $10 e-2$ | $9.7133 e-05$ | $1.0189 e-01$ | $2.3837 e-01$ |
|  | $10 e-5$ | $3.7580 e-05$ | $9.8341 e-04$ | $2.3096 e-03$ |
| 900 | $10 e-2$ | $4.0489 e-05$ | $1.0180 e-01$ | $2.3888 e-01$ |
|  | $10 e-5$ | $3.6766 e-05$ | $6.7920 e-04$ | $3.1624 e-03$ |

According to the results of tables below (3), (4), (5), (6), we can see that:

- Solutions cannot be reconstructed from partial knowledge by Tikhonov regularization or hybrid method even from 900 points of a priori data.
- From 25 points of a priori knowledge we can reconstruct PDE's solution with error $_{u}=2.2934 e-02$ by deep learning approach.
- With maximal a priori knowledge, DL-IP is more stable to noise $b$ than classical methods. With $b=10 e-3$, we get a solution with error $_{u}=2.7851 e-05$ by DL-IP on testing data, and with error $_{u}=2.0916 e-01$ by regularization approach.
- For classical methods we should use a percent of a priori knowledge calculated as $h \%=\frac{100 h}{S}$. Since our tests indicate that to get some precision the amount increases by increasing size of mesh.
- Solution reconstructed by hibrid approach from total data is better than the solution of Tikhonov regularization.


## 9. Conclusion

We have presented in this paper a new approach based on deep learning for the determination of parameters $P=\left(a, u_{0}, f\right)$ in linear parabolic equation from final data observations and partial a priori knowledge. Since the classical methods are not meshfree and require total a priori knowledge, our algorithm covers a great need in the domain of inverse problems.

Deep learning could become a valuable approach for solving inverse problems, which are important in physics, engineering, and finance. The PDE solution and its parameters can be approximated with deep neural networks, which are trained to satisfy the differential operator, initial condition, final observations, a priori knowledge, and boundary conditions. Our deep learning algorithm for solving inverse problems is mesh-free, which is key since meshes become infeasible in higher dimensions. Instead of forming a mesh, the neural network is trained on batches of randomly sampled time and space points. The algorithm can be easily modified to apply to hyperbolic, elliptic, and nonlinear equations. The algorithm remains the same for these other types of PDEs. However, the numerical performance of these other types of PDEs remains to be investigated. Developing deep learning algorithms that do not require any a priori knowledge or at least require final a priori knowledge would certainly be interesting. We leave this question for future work.
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