Machine learning approach for quantum non-Markovian noise classification

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In this paper, machine learning and artificial neural network models are proposed for quantum noise classification in quantum dynamics affected by external noise. For this purpose, we train and then validate support vector machine, multi-layer perceptron and recurrent neural network models with different complexity and accuracy, to solve supervised binary classification problems. As a result, we demonstrate the high efficacy of such tools in classifying noisy quantum dynamics using data sets collected from realizations of the quantum system dynamics. In addition, we show that for a successful classification one just needs to measure, in a sequence of discrete time instants, the probabilities that the analysed quantum system is in one of the allowed positions or energy levels. Several techniques, at both the theoretical and experimental side, have been developed for the inference of long-range two-qubit correlations. Moreover, Machine Learning (ML) models have been also adopted to study non-Markovian open quantum dynamics. In particular, in Ref. [31] a deep neural network approach is adopted to perform (at the theoretical level) noise regression of qubits immersed in their environment that entails different stationary, Gaussian noise spectra. In [31], deep neural networks are trained with time-dependent coherence decay curves used as input data.

In this paper, we exploit ML techniques to efficiently carry out quantum noise classification with high accuracy. The proposed methods are designed to distinguish between Independent and Identically Distributed (i.i.d.) noise sequences and noise samples originated by a non-trivial memory kernel, thus characterised by specific time-correlation parameters. It is worth reminding that, in the latter case, the dynamics of the stochastic quantum system (stochastic due to the presence of fluctuating parameters, e.g. in the Hamiltonian of the analyzed system as in [31]) turns out of being non-Markovian, in the sense that samples of its state in different time instants are correlated. This entails that the propagation of the system in subsequent time intervals is highly influenced by its previous states, even occurring in the early stages of the dynamics. This effect corresponds to a two-fold exchange of information between the system and the external sources, which has thus applications for quantum sensing.

To present our novel approach and demonstrate its efficacy in discriminating Markovian and non-Markovian
II. STOCHASTIC QUANTUM DYNAMICS

Let us introduce the general physical framework to which our ML methods will be applied. For this purpose, we consider a quantum particle that randomly moves on a complex graph $\mathcal{G}$ by following the quantum mechanics postulate. The complex graph is described by the pair $(\mathcal{N}, \mathcal{E})$, where $\mathcal{N}$ is the set of nodes or vertices while $\mathcal{E}$ is the set of edges linking the nodes. Each node of the graph $\mathcal{G}$ at discrete time instants, we will show that it is possible to discriminate accurately between different noise sources and identify the possible presence of time-correlations from observation of the quantum system dynamics.

To perform quantum noise classification, Support Vector Machines (SVMs), Multi-Layer Perceptrons (MLPs) and Recurrent Neural Networks (RNNs) [12, 45] are successfully trained on six data sets (each of them composed of 20,000 realisations) that have been properly generated to carry out binary classification of noisy quantum dynamics. Once trained, the proposed ML-models are able to reach a classification accuracy (defined by the number of correctly classified realisations over their total number) up to 97%. A pictorial representation of the proposed ML procedure is depicted in Fig. 1.

As other existing sensing techniques, the training of our ML models can be performed preliminary on synthetic data. Specifically, synthetic data are generated by solving a stochastic Schrödinger equation – modeling the noisy quantum dynamics we are analyzing – that exhibits at least one random parameter to be randomly sampled. Then, if useful, one shall improve subsequently the efficacy of the ML models once that further information on the experimental setup, as well as experimental data, are collected. This indeed allows to make the modelling more accurate by iterative optimizations in presence of known systematic errors. As a result, we have observed that both i.i.d. and correlated noise sources can be accurately discriminated by means of one single ML architecture. Moreover, our ML-based approach allows also for non-Markovian noise classification by processing only measurements of the diagonal elements (even called “populations”) of the density operator $\rho_t$ associated with the quantum system under investigation. Thus, no measurements of the off-diagonal elements of $\rho_t$, stemming from quantum coherence terms in a given basis of interest, might be required. For example, for the quantum particle case, this means that we just need to record, in discrete time instants, the probabilities (denoted as “occupation probabilities”) that the particle is in the positions (even part of them) identified by the nodes of the graph $\mathcal{G}$. These advantages constitute the novelty core of our paper, which is expected to find applications in experimental setups affected by stochastic noise sources as the ones in [9, 10, 12], and possibly in the available or coming quantum devices where a noise certification could be crucial before performing any task [14].

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is the set of links, denoted as $s \leftrightarrow \ell$, coupling pairs of
nodes, with $s, \ell = 1, \ldots, d$ and $d$ being the total number
of nodes. Each node is associated with a different parti-
cle position, while the links correspond to the possibility
that the particles jumps from a node to another. In partic-
ular, the links in $\mathcal{E}$ can be summarised in the adjacency
matrix $A_t$ (time-dependent operator in the more general
case), whose elements are given by

$$A_t(s, \ell) = \begin{cases} g_t \text{ if } s \leftrightarrow \ell \in \mathcal{E} \\ 0 \text{ if } s \leftrightarrow \ell \notin \mathcal{E}. \end{cases} (1)$$

In this way, we are implicitly assuming that all the links
are equally coupled with the same weight equal to $g_t$ that
is taken as a time-dependent parameter.

Here, the coupling $g_t$ is modelled as a stochastic pro-
cess defined by the collection of random variables $g \equiv
\{g_0, \ldots, g_{M-1}\}^T$, with $(\cdot)^T$ being the transposition
operation, in correspondence of the discrete time instants
$t_k, k = 0, \ldots, M - 1$. At each $t_k$, $g_k$ is sampled from
a specific probability distribution $\text{Prob}(g)$ and is assumed
to remain constant at the extracted value for the entire
time interval $[t_k, t_{k+1}]$. For simplicity, also the value $\Delta \equiv
\Delta_k = t_{k+1} - t_k$ is taken constant for any $k = 0, \ldots, M - 1$, and
the stochastic process $g_k$ is considered to take $D$ different
values $g^{(1)}, \ldots, g^{(D)}$ with probabilities $p_{g^{(1)}}, \ldots, p_{g^{(D)}}$. In this
way,

$$\text{Prob}(g) = \sum_{j=1}^{D} p_{g^{(j)}} \delta(g - g^{(j)}) (2)$$
is provided by a discrete probability distribution with $D$
values, with $\delta(\cdot)$ denoting the Kronecker delta.

If $g$ is provided by a collection of i.i.d. random vari-
ables sampled from the probability distribution $\text{Prob}(g)$,
then the noise sequence that affects the link strength $g$
is uncorrelated over time, and it is denoted as Marko-
vian. Conversely, in case the occurrence of the random
value $g^{(j)}, j = 1, \ldots, D$, at the discrete time instants
$t_k, k = 0, \ldots, M - 1$, depends on the sampling of $g$
at previous time instants, the noise sequence is time-
 correlated and the noise is denoted as non-Markovian
or as a coloured noise process. In this regard, notice that
the value of the parameters, which define the correlation
among different samples of noise in single time-sequences,
uniquely set the colour of the noise. Moreover, also ob-
serve that, known the multi-times distribution $\text{Prob}(g)$
defined over the discrete time instants $t_k$, one can com-
pute the noise auto-correlation function, whose Fourier
transform is by definition the power spectral density of
the noise process. In other terms, there is a one-to-one
mapping between the representations of the noise in the
time and frequency domains respectively. This entails
that noise sensing can be performed in one of the two
domain at best convenience. Moreover, this also moti-
vates the generality of the stochastic quantum model we
are here introducing that, indeed, can be applied to all
those problems concerning the transport of single parti-
cles within a network [46,48], but also to quantum sys-
tem dynamics influenced by the external environmental
as those in Refs. 2, 13, 13.

In our quantum dynamical model, we adopt as corre-
lation model the well-known formalism of time-
homogeneous discrete Markov chains [49]. The latter
can be graphically interpreted as a state-machine that
assign the conditional probability of “hopping” from each
possible value of $g$ to an adjacent one at consecutive
time instants. Each conditional probability is defined,
at any time $t$, by a transition matrix $T$ that is left or
right stochastic operator. Let us remind that discrete
Markov chains differ by a parameter $m$ named the order
of the chain. In a Markov chain of order $m$, future re-
alisations of the sampled random variable (e.g., our $g_t$
depend on the past $m$ realisations in previous time in-
stants. Here, we will consider $(m = 1)$-order discrete
Markov chains, namely correlated noise sequence charac-
terised by a single (1-step) transition matrix $T$ that we
aim to discriminate by means of properly-developed ML
techniques. This choice is simply dictated by our desire
to effectively illustrate the obtained results, and not by
intrinsic limitations of the methods we are going to pro-
pose. As an example, let us assume $m = 1$ and $D = 2$.
In this specific case, by taking the conditional probabil-
ites $p(g_k|g_{k-1})$ with $g_{k-1}$ equal to $g^{(1)}$ or $g^{(2)}$ for any $k$,
it holds that $p(g_k|g_{k-1})$ is equal to one of the elements
within the following transition matrix:

$$T = \begin{pmatrix} p(g_k = g^{(1)}|g_{k-1} = g^{(1)}) \quad p(g_k = g^{(2)}|g_{k-1} = g^{(1)}) \\ p(g_k = g^{(1)}|g_{k-1} = g^{(2)}) \quad p(g_k = g^{(2)}|g_{k-1} = g^{(2)}) \end{pmatrix} (3)$$

Thus, the stochastic realisations of $g$ in different time
instants are not correlated only if all the elements of $T$
are equal to $1/2$. In addition, we assume that all the nodes
of the graph $G$ have the same energy. Without loss of
generality, one is allowed to set such energy to zero, with
the result that the Hamiltonian $H_t$ of the quantum parti-
cle is identically equal to the adjacency matrix $A_t$, i.e.,
$H_t = A_t$ for any time instant $t$. Moreover, we consider
that the state of the particle, moving on a graph with
d nodes, is provided by the density operator $\rho_t$ that,
by definition, is an Hermitian, positive semi-definite, idem-
potent operator matrix with trace 1. By using the vec-
torisation operation vec$[\cdot]$, we convert $\rho_t$ into the column
vector

$$\lambda_t \equiv \text{vec}[\rho_t] = (\rho^{(1)}_t, \ldots, \rho^{(1)}_t, \rho^{(2)}_t, \ldots, \rho^{(d)}_t, \ldots, \rho^{(dd)}_t) \in \mathbb{C}^{d^2}$$

where $\rho^{(s, \ell)}_t$ denotes the $(s, \ell)$-element of $\rho_t$. The state
$\lambda_t$ is a vector of $d^2$ elements belonging to the space
of complex numbers. Since a quantum particle can live in a
superposition of positions, whereby also quantum coher-
ence plays an active role, $d$ elements of $\lambda_t$ corresponds
to the probabilities of measuring the particle in each of
the allowed positions, while the other elements are the so-called quantum coherence terms that identify interference patterns between the nodes of the graph. Thanks to the vectorisation of $\rho_t$, the ordinary differential equation, governing the dynamics of the particle, is recast in a linear differential equation for $\lambda_t$, i.e.,

$$\frac{\partial}{\partial t} \lambda_t = \mathcal{L}_t \lambda_t \iff \lambda_t = e^{\mathcal{L}_t} \lambda_0$$

$$\mathcal{L}_t \equiv -\frac{i}{\hbar} (\mathbb{I}_d \otimes A_t - A_t^T \otimes \mathbb{I}_d)$$

with $\otimes$ Kronecker product and $\hbar$ reduced Planck constant. By construction, $\mathcal{L}_t$ is a skew-Hermitian operator for any time instant $t$, i.e., $\mathcal{L}_t^\dagger + \mathcal{L}_t = 0 \ \forall t$.

### III. PROBLEM FORMULATION

Our aim is to identify the presence of noise sources acting on the coupling $g_t$ of the adjacency matrix $A_t$, and then discriminate among different noise probability distributions $\text{Prob}(g)$ and correlation parameters in the samples of the time-sequences $g$. Moreover, we also aim to evaluate if such tasks can be carried out by only measuring the population terms of the particle at the discrete time instants $t_k$, even by taking into account few runs of the quantum system dynamics. The population values are collected in the vectors $\mathcal{P}_{t_k} \in \mathbb{R}^d$ that have as many elements as the nodes of the graph. After each stochastic evolution of the quantum particle, $\mathcal{P}_{t_k}$ takes different values depending on the specific realisation of $g$. At the experimental level, the population distributions $\mathcal{P}_{t_k}$ can be obtained in multiple runs, by stopping the stochastic evolution of the system at each time $t_k$ (with $k = 1, \ldots, M$), then collecting the measurement records and restarting from the beginning the experimental routine.

#### A. Data set generation

For the generation of the data used to train the ML-models, we consider two variants of three different classification problems. Each sample of the data sets is created by first generating a random set of links $\mathcal{E}$ (random topology) for the graph $\mathcal{G}$, and then initialising the particle in a randomly chosen node of the graph. We set $M = 15$ as the number of evaluations (measurements) of the quantum particle dynamics, and $d = 40$ as the number of nodes of the graph $\mathcal{G}$. This means that $\mathcal{P}_{t_0}$ is a Kronecker delta centered in one of the 40 nodes, and the stochastic quantum dynamics is evolved for 15 steps for each simulated noise source of the generated data set. Here, it is worth noting that the choice of $M = 15$ is dictated by the fact that in recent experiments with monitored quantum systems, as for instance in Refs. [3, 50, 51], the number of intermediate quantum measurements does not exceed 10, and thus $M = 15$ is sufficiently large to represent actual physical setups. Instead, regarding taking $d = 40$, such a value is just able to generate a complex landscape for the particle dynamics and small enough to be numerically manageable. The total considered dynamical time $t_M$ is taken equal to $t_{15} = 1$ or $t_{15} = 0.1$ in dimensionless units, each of them corresponding to a specific variant.

Notice, indeed, that the values of the dynamical time $t_{15}$ are expressed consistently with the energy scale of the couplings $g_t$, whose random values $g^{(j)}$ in the data set generation belong to the set $\{1, 2, 3, 4, 5\}$, such that $\hbar$ can be reliably set to 1 as usual. All the probability distributions $\mathcal{P}_{t_k}$ for $k = 0, \ldots, 15$ are stored together with the attached label that indicates the associated type of noise.

For each of the two variants of classification problems, we generate three different balanced data sets of 20 000 samples. The first data set, which we call IID, is suitable for a supervised binary classification task that discriminates between two different i.i.d. noisy quantum dynamics, where the noise sources have the same support but different probability distribution $\text{Prob}(g)$.

The second data set, named as NM, concerns the classification of two different coloured noisy quantum dynamics with the noise sources again having the same support but different $\text{Prob}(g)$ (the same ones as in the data set IID) and a transition matrix $T$.

Finally, the third data set, called VS, is created for the classification between stochastic quantum dynamics affected respectively by an i.i.d. and a coloured noise with same support and $\text{Prob}(g)$.

Note that choosing graphs with random links allows to increase the statistical variability of the input data, with the result that the ML algorithms learn to classify noise sources independently of the graph topology. The aim, indeed, is to prevent that the ML-models rely only on features specific to a small class of topologies. Moreover, taking random initial distributions $\mathcal{P}_{t_0}$ allows to increase the robustness of the ML methods, making them less likely to overfit on the synthetic data set.

As it will be explained in the following, some ML-models that we are going to introduce will use as input only the last distribution $\mathcal{P}_{t_{15}}$, while other ML-models will take all the $\mathcal{P}_{t_k}$ for any $t_k$. Moreover, each data set is balanced split in a training set of 12 000 samples, a validation set of 4 000 samples, and a test set of 4 000 samples.

In Tables II and III we plot the occupation probabilities $\mathcal{P}_{t_k}$ (just for the IID case for the sake of an easier presentation), being here interested in looking for the difference between choosing $t_{15} = 0.1$ or 1, which identify the two different variants of the generated data set. In this regard, it is worth noting that the duration $t_{15} = 0.1$ (in dimensionless units) of the quantum system dynamics, as in the example in Table II, is the minimal one to observe the diffusion of the system’s population outside the node on which has been initialised. However, as it will be verified by our experiments and explained later, with this choice one has that, by taking $t_{15} = 0.1$, the classification problem results quite straightforward. Indeed, just ba-
Table I. Example of a part of $\mathcal{P}_{t_k}$ for all the discrete time instants $t_k$ for a noisy quantum dynamics affected by i.i.d. noise sources and $t_{15} = 0.1$ (in dimensionless units). In the Table, $\mathcal{P}_{t_k}^{(s)}$ denotes the $s$-th element of the vector $\mathcal{P}_{t_k}$ for any $t_k$, $k = 0, \ldots, 15$.

| $\mathcal{P}_{t_k}^{(35)}$ | $\mathcal{P}_{t_k}^{(36)}$ | $\mathcal{P}_{t_k}^{(37)}$ | $\mathcal{P}_{t_k}^{(38)}$ | $\mathcal{P}_{t_k}^{(39)}$ | $\mathcal{P}_{t_k}^{(40)}$ |
|--------------------------|--------------------------|--------------------------|--------------------------|--------------------------|--------------------------|
| $t_0$                    | 0.00                     | 0.00                     | 0.00                     | 1.00                     | 0.00                     |
| $t_1$                    | 0.00                     | 0.00                     | 0.00                     | 0.99                     | 0.00                     |
| $t_2$                    | 0.00                     | 0.00                     | 0.00                     | 0.93                     | 0.00                     |
| $t_3$                    | 0.00                     | 0.01                     | 0.01                     | 0.85                     | 0.01                     |
| $t_4$                    | 0.01                     | 0.01                     | 0.01                     | 0.78                     | 0.01                     |
| $t_5$                    | 0.01                     | 0.01                     | 0.01                     | 0.69                     | 0.01                     |
| $t_6$                    | 0.01                     | 0.02                     | 0.01                     | 0.63                     | 0.00                     |
| $t_7$                    | 0.01                     | 0.02                     | 0.01                     | 0.57                     | 0.00                     |
| $t_8$                    | 0.01                     | 0.02                     | 0.01                     | 0.52                     | 0.00                     |
| $t_9$                    | 0.02                     | 0.02                     | 0.01                     | 0.45                     | 0.00                     |
| $t_{10}$                 | 0.02                     | 0.02                     | 0.02                     | 0.37                     | 0.01                     |
| $t_{11}$                 | 0.01                     | 0.02                     | 0.02                     | 0.29                     | 0.01                     |
| $t_{12}$                 | 0.01                     | 0.01                     | 0.03                     | 0.20                     | 0.02                     |
| $t_{13}$                 | 0.01                     | 0.04                     | 0.04                     | 0.14                     | 0.01                     |
| $t_{14}$                 | 0.01                     | 0.04                     | 0.05                     | 0.08                     | 0.01                     |
| $t_{15}$                 | 0.01                     | 0.04                     | 0.05                     | 0.06                     | 0.01                     |

Table II. Example of a part of $\mathcal{P}_{t_k}$ for all the discrete time instants $t_k$ for a noisy quantum dynamics affected by i.i.d. noise sources and $t_{15} = 1$ (in dimensionless units). Again, $\mathcal{P}_{t_k}^{(s)}$ denotes the $s$-th element of the vector $\mathcal{P}_{t_k}$ for any $t_k$, $k = 0, \ldots, 15$. The topology and the initial state, for this example, are the same as those in Table I.

| $\mathcal{P}_{t_k}^{(35)}$ | $\mathcal{P}_{t_k}^{(36)}$ | $\mathcal{P}_{t_k}^{(37)}$ | $\mathcal{P}_{t_k}^{(38)}$ | $\mathcal{P}_{t_k}^{(39)}$ | $\mathcal{P}_{t_k}^{(40)}$ |
|--------------------------|--------------------------|--------------------------|--------------------------|--------------------------|--------------------------|
| $t_0$                    | 0.00                     | 0.00                     | 0.00                     | 1.00                     | 0.00                     |
| $t_1$                    | 0.02                     | 0.02                     | 0.01                     | 0.45                     | 0.00                     |
| $t_2$                    | 0.01                     | 0.01                     | 0.02                     | 0.03                     | 0.05                     |
| $t_3$                    | 0.00                     | 0.00                     | 0.00                     | 0.13                     | 0.02                     |
| $t_4$                    | 0.02                     | 0.01                     | 0.01                     | 0.12                     | 0.02                     |
| $t_5$                    | 0.01                     | 0.01                     | 0.03                     | 0.06                     | 0.01                     |
| $t_6$                    | 0.01                     | 0.01                     | 0.01                     | 0.01                     | 0.00                     |
| $t_7$                    | 0.04                     | 0.03                     | 0.01                     | 0.11                     | 0.00                     |
| $t_8$                    | 0.04                     | 0.00                     | 0.03                     | 0.11                     | 0.03                     |
| $t_9$                    | 0.03                     | 0.00                     | 0.03                     | 0.01                     | 0.10                     |
| $t_{10}$                 | 0.05                     | 0.01                     | 0.04                     | 0.08                     | 0.04                     |
| $t_{11}$                 | 0.01                     | 0.03                     | 0.02                     | 0.00                     | 0.08                     |
| $t_{12}$                 | 0.00                     | 0.05                     | 0.02                     | 0.04                     | 0.00                     |
| $t_{13}$                 | 0.01                     | 0.03                     | 0.00                     | 0.02                     | 0.05                     |
| $t_{14}$                 | 0.00                     | 0.00                     | 0.00                     | 0.01                     | 0.12                     |
| $t_{15}$                 | 0.00                     | 0.00                     | 0.01                     | 0.04                     | 0.10                     |

Table II shows that the trained ML-models that are only trained on $\mathcal{P}_{t_{15}}$ (thus, only on the final distribution $\mathcal{P}$) are able to correctly classify between two noisy quantum dynamics. Therefore, it was more interesting to increase the value of $t_{15}$ up to $t_{15} = 1$ (in dimensionless units). As in the example of Table III it leads to more complex data sets, and only deep learning models, designed to read all the $\mathcal{P}_{t_k}$, can classify the generated noisy quantum dynamics. According to preliminary simulations, with a value of $t_{15} \gg 1$ our trained ML-models do not reach more than 60% of accuracy for the test set.

As final remark, note that the current synthetic data set is build assuming perfect measurement statistics, as it was obtained from a large enough number of repetitions of the noisy quantum dynamics. Hence, to better adapt the synthetic data set to real data, one should simulate experimental case in which the measurement statistics are estimated from a finite number of dynamics realizations (i.e., measurement shots).

IV. MACHINLEARNING MODELS

In this section we present and compare the employed classical machine learning algorithms and more recent deep learning approaches to perform binary supervised classification tasks. We recall that, in our case with $\mathcal{P}_{t_k}$ as input, we aim to classify between:

(i) Two different probability distributions

Prob$(g)$ – specifically, $p_g^{(1)}, \ldots, p_g^{(5)} = (0.0124, 0.04236, 0.0820, 0.2398, 0.6234)$ and

$= (0.1782, 0.1865, 0.2, 0.2107, 0.2245)$ – both associated with i.i.d. noise sources.

(ii) Two different Prob$(g)$ (the same as (i)) and different values of the correlation parameters – identified by transition matrices $T$ as explained in Sec III – for coloured (thus, non-Markovian) noise processes.

(iii) An i.i.d. and a coloured noise process with the same support $g$ and distribution $p_g^{(1)}, \ldots, p_g^{(5)} = (0.0124, 0.04236, 0.0820, 0.2398, 0.6234)$ that thus differ for the presence of non-zero correlation parameters.

The values of both Prob$(g)$ and the transition matrices $T$, used in our numerical simulations, are chosen randomly.

A generic binary data set in input to ML-models is usually represented by a set of $n$ points $x_q \in \mathbb{R}^p$, with $q = 1, \ldots, n$, each of them living in the $p$-dimensional space of the features. A feature is a distinctive attribute of each element of the data set. Each point $x_q$ is associated with one of two different classes with binary labels $y_q \in \{-1, 1\}$, with $q = 1, \ldots, n$, depending on the specific classification problem that we are solving.

Support Vector Machine (SVM) \cite{33} is a classical ML model that is used for classification tasks. The SVM training consists in finding the hyperplane that separates the elements $x_q$ in two groups: one with the label $y_q = 1$ and the other with $y_q = -1$. The final hyperplane, solution of the classification, is the one having the maximum geometrical distance from the two parallel hyperplanes.
that are defined by the subsets of $\mathbf{x}_\tau$ called the support sets. When the data is not linearly separable, the kernel trick allows to increase the dimension of the features space in a way that the data becomes linearly separable in the new space.

There are several classification problems (as for example the ImageNet Large Scale Visual Recognition Challenge [52], employing millions of images with hundreds of categories) that are solved through SVM but without a low enough residual classification error. For this reason, to improve the performance in solving classification problems, Artificial Neural Networks (ANNs) have recently been (re-)introduced as more-performing tools, and since 2012 have been extensively used [44, 52–54]. In general a Multi-Layer Perceptron (MLP), which is the basic form of ANN, is defined by

$$\hat{y} = f(\mathbf{x}; \theta, \xi),$$

where $\mathbf{x}$ is the input, $\hat{y}$ denotes the predicted output, and $f$ is a non-linear function composed of several layers of artificial neurons parametrized with weights $\theta$ and with a structure defined by the hyperparameters $\xi$. The training of the parameters $\theta$ is performed minimizing the error done by the model on the training set, by means of Stochastic Gradient Descent (SGD) techniques. Instead, the hyperparameters $\xi$ are tuned using hyperparameters optimization methods that minimize the error over the validation set.

A Recurrent Neural Network (RNN) is an ANN specialised for sequence processing when the data set is expressed as

$$S = \{\{\mathbf{x}_1, \ldots, \mathbf{x}_\tau\}, \mathbf{y}_1, \ldots, \{\mathbf{x}_1, \ldots, \mathbf{x}_\tau\}, \mathbf{y}_n\},$$

where $\tau_r$ defines the number of elements of the $r$-th sequence. RNNs can be used in tasks regarding Natural Languages [42–44, 52]. One of the most common techniques used to develop RNN are Long Short Term Memory (LSTM) [67] and Gated Recurrent Unit (GRU) [68]. Both of them adopt a gated mechanism to dynamically decide (i) how much of the input elements has to be processed to update the hidden state of the RNN, and (ii) how much of the hidden state has to be used to produce the intermediate outputs. The classification task is performed by a MLP applied to the output of the RNN. By taking into account the definition in Eq. (6), the output can be represented by the last representation $\mathbf{h}_r$ or, more in general, by the aggregation of all the $\mathbf{h}_1, \ldots, \mathbf{h}_r$. Such aggregation is performed with the help of the so-called attention mechanisms [69–72] or the simpler max pooling. Aggregation with max pooling has been adopted in NLP tasks, and, under certain conditions, has given interpretability to the RNN models [52, 60].

For an exhaustive explanation of the ML models employed in this paper, refer to Appendix.

V. RESULTS

In our work, we consider two SVM models as baseline. The first one is denoted $\mathbf{m}$-SVM-single and uses as input only the final probability distribution $P_{t15}$ (the prefix $m$-stands for “model”, to avoid confusion with the algorithm name; the suffix -single means that it is based only on $P_{t15}$). Instead, the second one, which we call as $\mathbf{m}$-SVM, uses the set of all the $P_{i}$, with $k = 0, \ldots, 15$. For both of them, we try the following kernels to increase the dimension of the feature-space that makes linearly separable the data-set: linear, polynomial with degree 2, 3, and Radial Basis Function (RBF).

Then, we denote with $\mathbf{m}$-MLP-single a MLP in the form of Eq. (4), with $\mathbf{x} \equiv P_{t15}$ and $\mathbf{y} \equiv (0, 1)$ or (1, 0) to identify the two noisy quantum dynamics that we aim to classify. Differently, $\mathbf{m}$-MLP takes as input the set of all $P_{i}$.

Moreover, $\mathbf{m}$-GRU and $\mathbf{m}$-LSTM are unidirectional RNNs that employ the final hidden representation (refer to Eqs. (21) and (22) in Appendix for more details). They are implemented by exploiting the GRU and LSTM methods, respectively. The input to the models is $\mathbf{x}_{i+1} \equiv P_{i}$, with $i = 0, \ldots, 15$, while the output $\mathbf{y} \equiv (0, 1)$ or (1, 0) as before. Besides, $\mathbf{m}$-biGRU and $\mathbf{m}$-biLSTM are the bidirectional versions of $\mathbf{m}$-GRU and $\mathbf{m}$-LSTM, while $\mathbf{m}$-biGRU-att and $\mathbf{m}$-biGRU-max are as $\mathbf{m}$-biGRU but in addition, respectively, with the attention mechanism and the max pooling (respectively, Eqs. (20) and (27) in Appendix) as forms of aggregation. Similarly, $\mathbf{m}$-biLSTM-att and $\mathbf{m}$-biLSTM-max are, respectively, the attentive and max pooling equivalents of $\mathbf{m}$-biLSTM.

In Table III for each model we report the best classification accuracy that is computed on the predictions performed over the test set. More formally, we define the prediction set

$$\Gamma \equiv \{(\mathbf{y}_1, \hat{\mathbf{y}}_1), \ldots, (\mathbf{y}_n, \hat{\mathbf{y}}_n)\}$$

where $\mathbf{y}_1, \ldots, \mathbf{y}_n$, taken from the data set, denote the true noise sources affecting the quantum system dynamics, and $\hat{\mathbf{y}}_1, \ldots, \hat{\mathbf{y}}_n$ the corresponding predictions of the
Table III. Percent accuracy $\gamma$ (calculated on the test set) of the ML-models trained in the tasks of binary classification of noisy quantum dynamics with: (i) Two different i.i.d. noise sources (IID); (ii) two different coloured noise processes (NM) leading to non-Markovian dynamics; and (iii) one i.i.d. vs one coloured noise sources (VS). In this regard, let us recall that the coloured noise processes addressed in this paper are such that the probability distributions $P_{t_k}$ depend both on $\text{Prob}(g)$ and $I$-step transition matrix $T$.

In the first three columns of the table, the total duration of the dynamics is equal to $t_{15} = 0.1$, while in the last three is $t_{15} = 1$. The first two rows of the table report the results of the ML-models that use as input only $P_{t_{15}}$, while the models of the other rows take as input all the probability distributions $P_{t_k}$ for $k = 0, \ldots, 15$. The highest values of the accuracy have been underlined, and a color gradient (from blue to bright red) highlights the difference in their values.

| $P_{t_{15}}$ | m-SVM-single | m-MLP-single | m-LSTM | m-biLSTM | m-biGRU | m-biGRU-att | m-biGRU-max | m-biLSTM-max |
|--------------|---------------|---------------|---------|-----------|----------|-------------|-------------|--------------|
| IID NM VS    | 97.0 82.3 96.5 | 96.9 80.7 96.6 | 96.5 91.5 96.7 | 96.8 90.4 96.4 | 96.6 92.2 96.6 | 96.9 87.9 96.3 | 96.6 92.6 96.6 | 96.6 91.4 96.3 |
| $I_{15} = 1$ | 96.3 51.2 49.5 | 90.4 60.4 49.5 | 90.5 50.7 50.2 | 90.5 73.3 88.2 | 90.5 88.6 86.3 | 90.9 74.6 90.6 | 90.9 76.1 90.4 | 91.4 74.9 89.0 |

Table III. Percent accuracy $\gamma$ (calculated on the test set) of the ML-models trained in the tasks of binary classification of noisy quantum dynamics with: (i) Two different i.i.d. noise sources (IID); (ii) two different coloured noise processes (NM) leading to non-Markovian dynamics; and (iii) one i.i.d. vs one coloured noise sources (VS). In this regard, let us recall that the coloured noise processes addressed in this paper are such that the probability distributions $P_{t_k}$ depend both on $\text{Prob}(g)$ and $I$-step transition matrix $T$.

In the first three columns of the table, the total duration of the dynamics is equal to $t_{15} = 0.1$, while in the last three is $t_{15} = 1$. The first two rows of the table report the results of the ML-models that use as input only $P_{t_{15}}$, while the models of the other rows take as input all the probability distributions $P_{t_k}$ for $k = 0, \ldots, 15$. The highest values of the accuracy have been underlined, and a color gradient (from blue to bright red) highlights the difference in their values.

| $P_{t_{15}}$ | m-SVM-single | m-MLP-single | m-LSTM | m-biLSTM | m-biGRU | m-biGRU-att | m-biGRU-max | m-biLSTM-max |
|--------------|---------------|---------------|---------|-----------|----------|-------------|-------------|--------------|
| IID NM VS    | 97.0 82.3 96.5 | 96.9 80.7 96.6 | 96.5 91.5 96.7 | 96.8 90.4 96.4 | 96.6 92.2 96.6 | 96.9 87.9 96.3 | 96.6 92.6 96.6 | 96.6 91.4 96.3 |
| $I_{15} = 1$ | 96.3 51.2 49.5 | 90.4 60.4 49.5 | 90.5 50.7 50.2 | 90.5 73.3 88.2 | 90.5 88.6 86.3 | 90.9 74.6 90.6 | 90.9 76.1 90.4 | 91.4 74.9 89.0 |

ML-model. Hence, the (percent) accuracy $\gamma$, function of $\Gamma$, is provided by

$$\gamma(\Gamma) = \frac{100}{n} \sum_{i=1}^{n} \mathbb{I} \left\{ \arg \max_{j=1,2} y_{i}^{(j)} = \arg \max_{j=1,2} y_{i}^{(j)} \right\}, \quad (7)$$

where

$$\mathbb{I}(c) = \begin{cases} 1, & \text{if } c \text{ is true,} \\ 0, & \text{otherwise.} \end{cases} \quad (8)$$

is the so-called indicator function. The accuracy $\gamma$ defines the correctness of the model and can be used as a metric to identify which solution is better. In detail, for binary classification problems, as in our case, if $\gamma \approx 50$ the classification is equivalent to perform a random guess, thus the model does not work. Instead, when $\gamma \approx 100$ the model perfectly classifies all the elements of the test set, thus it is a nearly ideal classifier.

From the table, one can first observe that, by dealing with a total duration of the dynamics (in dimensionless units, by rescaling as the inverse of the couplings $g_t$) equal to $t_{15} = 0.1$, we can reach the 97% and 96.6% of accuracy for the classification tasks IID and VS via an SVM using as input only the distribution $P_{t_{15}}$. Instead, the task NM is more difficult: 82.3% of accuracy is achieved by SVMs applied just on $P_{t_{15}}$. MLP does not provide better results. In this case (NM tasks), to obtain an accuracy over 90%, one can resort to RNNs taking as inputs all the $P_{t_k}$ for $k = 0, \ldots, 15$.

Conversely, for a longer dynamics, i.e., with $t_{15} = 1$, we notice that using only $P_{t_{15}}$ all the three classification tasks are not solved neither with SVM or MLP. Indeed, the accuracy $\gamma$ is always around 50% and the models basically perform random guesses. The accuracy is increased by means of an SVM or an MLP based on all $P_{t_k}$, with $k = 0, \ldots, 15$ as input. However, to get over 90% of accuracy on the tasks IID and VS, we need to employ RNNs. The task NM with $t_{15} = 1$ is the most difficult among the analysed ones, and just 76.1% of accuracy is obtained using RNNs. It is worth noticing that, for the tasks with $t_{15} = 1$, we have empirically observed that the models adopting GRU perform better with respect to the ones that employ LSTM. Moreover, setting the bidirectionality in the RNNs allows slight improved accuracy, as well as the use of max pooling in aggregation. Instead, the attention mechanism does not seem to be beneficial for these tasks.

Among the proposed solutions, the more-performing is m-biGRU-max that is realised by a bidirectional RNN with GRU and max pooling aggregation. However, from our numerical simulations, we have observed that, independently on the employed ML model, the value of the total dynamical time as well as $M$ and $\Delta$ (see also paragraph [V.A] emerge to be crucial for quantum noise classification. Specifically, by taking a quantum dynamics with a short enough duration, also SVMs are able to classify quantum noise sources with very high accuracy. With short enough dynamics we mean short with respect to the time needed to the particle in escaping from the initial node of the graph, which in our case is around $t_{15} = 0.1$. Instead, with $t_{15}$ around 1 only RNNs provide better results, and for $t_{15} \ll 1$ none of the proposed ML-techniques solves quantum noise classification problems (these results have not been reported in Table III for the sake of better presentation). It is also worth stressing that, if the duration of the quantum dynamics is $t_{15} = 0.1$, ML-models efficiently classify quantum noise sources by only processing the last measured distribution $P_{t_{15}}$. These findings can be relevant for effective implementation (also at the experimental level), since the training and tuning of SVM is orders of magnitude faster with respect to ANNs (e.g., around minutes vs hours or even days depending on the model and provided that a GPU is used). The reason to that has to be found in the more complex structure of the ANNs than SVMs.

A. Scaling of the classification accuracy

Let us now investigate the scaling of the classification accuracy $\gamma$, as a function of both the interval $\Delta$ between two consecutive transitions for $g$ and the number $M$ of discrete time instants. Notice that $\Delta$ and $M$ are related
to the total dynamical time $t_M$, since $t_M \equiv M\Delta$.

A possible explanation of the differences observed between the three previously-analysed scenarios, i.e., $t_{15} = 0.1$, $t_{15} = 1$ and $t_{15} \gg 1$ (in dimensionless units), could be that the information on both the noise source and the initial quantum state is lost during the evolution of the system. For such aspect, not only the total dynamical time $t_{15}$ could play a role, but also the time interval $\Delta \equiv t_1 - t_0 \equiv \cdots \equiv t_M - t_{M-1}$. In fact, it is reasonable to conjecture that a ML-model, able to correctly classify our noisy quantum dynamics with $t_{15} = 1$ (thus $M = 15$), can also work with $t_{M'} \gg 1$ for $M' > 15$ and $\Delta' = \Delta$ where $\Delta' \equiv t_1 - t_0 \equiv \cdots \equiv t_{M'} - t_{M'-1}$. In this way, the sequence $\mathcal{P}_{t_1}, \ldots, \mathcal{P}_{t_{15}}$ is contained in $\mathcal{P}_{t_1}, \ldots, \mathcal{P}_{t_{M'}}$. In other terms, we conjecture that the classification problem can be solved even for longer noisy quantum dynamics, but provided that $\Delta$ remains small.

To gain evidence on this conjecture, we have performed two additional experiments. Starting from the task IID with $t_{15} = 1$ and m-biGRU-max as baseline (accuracy 91.8%), the same model (optimised in the same hyperparameters space) is trained on two new data sets. In both data sets, $t_M = 2$ with $M$ equal to 15 for the first data set and 30 for the second one. Thus, in the former $\Delta' > \Delta$ with $\Delta$ time interval of the original data set, while in the latter $\Delta' = \Delta$. The first experiment ($\Delta' > \Delta$) provides a classification accuracy of 81.1%, contrarily to the results from the second experiment ($\Delta' = \Delta$), where a better accuracy of 96.3% is achieved. We thus observe that, by taking $\Delta' = \Delta$ and the same ML-model, the classification problem can be solved with an higher accuracy, but at the price of a longer training time. Indeed, in this case, the length of each sample of the data set is twice the original one.

In another experiment, whose results are shown in Figure 2, we vary $M$ by keeping the total evolution time equal to $t_M = 1$. Such tests use as baseline the model m-biGRU-max applied on the most difficult task of Table III, i.e., NM with $t_{15} = 1$. As a result, the achieved classification accuracy is directly proportional to $M$ and, thus, inversely proportional to the value of the time interval $\Delta$. Indeed, by taking $t_M$ fixed and reducing $\Delta$, the classification accuracy of the same model can be enhanced. Specifically, it is possible to obtain more than 90% of accuracy also for the task NM with a total dynamical time equal to 1, at the price of a longer training time as the length of the sequences increases.

B. Quantum advantages

Here, we address the following question: Could the proposed ML techniques be applied for the inference of noise sources affecting the dynamics of classical systems, e.g., Langevin equations? Probably yes, but we expect that their application to quantum systems, maybe surprisingly, can be more effective than on classical systems. Both classical (non-periodic) dissipative dynamics and stochastic quantum dynamics (stochastic due to the presence of an external environment, or noise sources as in our case) can asymptotically tend to a fixed-point, whereby the information on the initial state is lost. This means that the states of the system used for this noise classification tend to become indistinguishable as time increases. Classically, this can happen due to energy dissipation introduced by damping terms. Instead, quantum-mechanically, a dynamical fixed point can be reached due to decoherence that makes vanishing, at least on average, all quantum coherence terms. Thus, once the transient of the evolution is elapsed, the evaluation of the final state of the system does not bring information neither on the initial state nor on the initial dynamics bringing the system to the asymptotic fixed-point. In our case, we have observed that, by using only $\mathcal{P}_{t_{15}}$ with long total dynamical time, the accuracy of all the classification tasks is always around 50% both for SVM and MLP. Consequently, if one aims to infer/reconstruct the value of parameters, signals or operators that influence the system dynamics by measuring its evolution, the most appropriate time window is during the transient. In this regard, a quantum dynamic, until it is nearly close of being unitary, is able to explore different configurations thanks to linearity and the quantum superposition principle. Conversely, classical dynamics, not being able to propagate superpositions of their trajectories, cannot provide per time unit the same amount of information on the quantity to be inferred.

In conclusion, the application of the proposed methods is expected to be more accurate if applied to quan-

![Figure 2. Percent classification accuracy $\gamma$ vs. $M$. It refers to the test set associated with the model m-biGRU-max for the task NM, where the value of the total evolution time is fixed to $t_M = 1$. It is worth noting that the first point of the figure corresponds to the value in Table III obtained for $t_{15} = 1$.](image-url)
tum systems than classical ones, but during the transient of its dynamical evolution when quantum effects are still predominant and the distance among the state and the fixed point is not negligible. This conjecture will be properly discussed, and possibly numerically proved, in a forthcoming paper.

C. Proposal for application to quantum computers

Our techniques are expected to be successfully employed for the classification of Markovian and non-Markovian noise sources in commercial quantum devices, as for example the Q-IBM® or Rigetti®. In fact, the latter, as other Noisy Intermediate-Scale Quantum prototypes [2], are unavoidably affected by the external environment that entails random errors. Recently, in Refs. [22, 72], it has been shown that it is possible to discriminate different quantum computers by looking at the (unknown) noise fingerprints that characterize each device. Thus, what the ML techniques – presented here – could provide as an added value is to evaluate whether such noise fingerprints are time-correlated, and possibly how much the time-correlation is non-Markovian. For such experimental noise benchmarking, as in [22], it could be convenient to fix the connections among quantum gates (i.e., the underlying topology), and then consider more realizations of the implemented quantum dynamics affected by noise. Then, the ML techniques could be implemented according to the following two steps:

1) Discriminate if and how much the measurement statistics (provided by the distributions \( P_k \) with \( k = 1, \ldots, M \)), which have been measured on the real quantum devices, differ from the corresponding theoretical predictions. Such a difference, here on denoted as \( D \), between theoretical and experimental data returns an effective prediction of the presence of noise on the machines.

2) Conditionally to step 1), classify with ML models the presence or the absence of functional relations \( \mathcal{F} \) that link the difference distributions \( D_{k_{1}} \) in correspondence of the time instants \( t_{k} \). If two consecutive instances of \( D \) at times \( t_{k-1} \) and \( t_{k} \) are no functionally related, then the noise is originated from a i.i.d. stochastic process. If, instead, there exist a functional \( \mathcal{F}_{t_{k-1}, t_{k}} = f(D_{t_{k-1}}, D_{t_{k}}) \) that links together in a non-trivial way \( D_{t_{k-1}} \) and \( D_{t_{k}} \), then the noise would come from a Markovian process. Finally, the noise process would be non-Markovian for functional relations \( \mathcal{F}_{t_{k-n}, t_{k}} \) defined over multi times, with \( n > 1 \).

The functional relations \( \mathcal{F} \) might be provided by non-linear functions \( f \) (as for example a polynomial development) parametrized by a given set of coefficients. Accordingly, the discrimination of functional relations among \( D_{t_{k}} \) would be equivalent to classify different sets of coefficients that define \( \mathcal{F} \).

To conclude, according to our proposal that will be tested in a forthcoming paper, time-correlations in the noisy samples of the distributions \( P_k \), with \( k = 1, \ldots, M \), can be determined by classifying functional relations \( \mathcal{F} \) linking the difference distributions \( D \), obtained by comparing the theoretical and measured values of \( P \) for a set of time instants. This is equivalent to discriminate coloured noise processes originated by different discrete Markov chain with non-zero transition matrices \( T \). Of course, also in this case, each time a projective measurement is performed and the resulting outcome recorded, the implemented quantum circuit shall be executed from the beginning.

VI. CONCLUSIONS

In this paper, we have addressed quantum non-Markovian noise classification problems by means of deep learning techniques. In particular, the use of RNN – developed for sequence processing – is motivated by the fact that we deal with time-ordered sequences of data. Even without resorting to external driving that may hinder detection tasks, we managed to classify with high accuracy stochastic quantum dynamics characterized by random parameters sampled from different probability distributions, associated with i.i.d. (Markovian) and coloured (non-Markovian) noise processes. For such a purpose, several ML models have been tested; in this regard, refer to Table [3] for a summary of the results in term of the classification accuracy.

Among the proposed solutions, the more-performing is **m-biGRU-max** that is realised by a bidirectional RNN with GRU and max pooling aggregation. In fact, recurrent neural networks are particularly suitable to accomplish temporal machine-learning tasks thanks to their capability to generate internal temporal dynamics based on feedback connections. However, independently on the employed ML model, different accuracy values are achieved depending on the values of \( M \), \( \Delta \) and the total dynamical time. The way our ML techniques rely on the parameters of the quantum model has been addressed in the paragraph [3]

Overall, all our numerical results have shown that it is easier to classify between two different noisy quantum dynamics both affected by i.i.d. noise sources or by i.i.d. and coloured noise processes than between two noisy quantum dynamics subjected to coloured noise. Again it confirms the relevant role played by time-correlations and how the latter highly influence the value of the classification accuracy. Furthermore, we also expect that the same ML techniques that we have exploited in this work could be successfully applied to classify among coloured quantum noise with \( q \)-step transition matrices \( T_{t_{i}, t_{i-q}} \) with \( q > 1 \).
A. Outlooks

As outlook, we plan to test the ML-models employed in this paper on reconfigurable experimental platforms as the ones in Refs. [77, 78], even affected by multiple noise sources. Moreover, we also aim to adapt our ML methods (and especially ANNs) to reconstruct quantum noise processes with time-correlation as key feature in the context of regression task instead of classification. Indeed, our proposal is to provide accurate estimates of both the probability distribution $\text{Prob}(g)$ and the transition matrix $T$, and the analysis would be extended for the prediction of spatially-correlated noise sources. In this way, ML approaches would represent a very promising, and possibly more accurate, alternatives to other noise-sensing techniques, e.g., those recently discussed in Refs. [77, 78].

A well-known problem in ML is the generalization to data shift. A model that is trained on a data set sampled from a specific data distribution will work correctly only with data sampled from the same distribution. In this paper, we used only synthetic data to evaluate the correctness of the training process and the ML techniques. Thus, in order to validate this approach to real data, we should first collect them. This is out of the scope of the current work, but, as a remark, we can delineate three possible ways to build a real experimental data set. The first strategy is to acquire information a-priori on noise sources affecting the quantum system of interest in some experimental contexts by means of standard spectroscopy techniques, so that we can train the proposed ML models to discriminate between unseen classes of noise. In this way, the initial effort in building a training data set that also contains experimental data is counterbalanced by the possibility to predict noise features by means of faster classification tasks. The second strategy, which has been employed in [72], is to collect experimental data that comes from different noisy measurement statistics whose noise processes are not necessarily known. The ML models, then, are trained to classify (unknown) noise sources in distinct unseen sets of measurements. Finally, the third strategy, which is aimed to reduce the effort in building an informative experimental data set, is to train the ML models first on synthetic data and then to fine tune the training on a smaller further data set with only experimental data. In such a case, it is beneficial to adopt a synthetic data set that closely adapts to the real experimental setup. For instance, a simulated extra training data set that closely adapts to the experimental data. In such a case, it is beneficial to train the ML models first on synthetic data and then to fine tune the training on a smaller further data set with synthetic data.

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DATA AND CODE AVAILABILITY

The source codes for the generation of the data sets and the ML experiments are available on GitHub at the following address:

https://github.com/trianam/quantumNoiseClassification

APPENDIX: DETAILS ON EMPLOYED ML MODELS

In this section, aiming at addressing also an audience not necessarily expert in ML, we describe more in detail the ML models used in our tasks.

Support Vector Machine

SVM is a generalisation of the Support Vector Classifier (SVC) that, in turn, is an improved version of the Maximal Margin Classifier (MMC) [33]. MMCs aim at finding the hyperplane separating the two aforementioned classes of points, such that the distance between the hyperplane and the nearest points of the classes (commonly denoted as margin) is maximised. If the points of the data set are not linearly separable, then the value of the margin is negative. In such a case, the MMCs cannot be adopted. SVCs increase the performance of MMCs, by allowing some points of the data set, called slack variables, to be in the opposite part of the hyperplane with respect to the others of the belonging class. If the data set exhibits a non-linear bound between the two classes of points, SVCs are not able to correctly separate them, albeit the method returns a solution. Finally, SVMs extend the capabilities of SVCs by increasing the number of dimensions of the feature-space, such that in the new space the data set becomes linearly separable.

Multi-Layer Perceptron

The basic model of the ANNs is the MLP. The latter is composed of a variable number of fully connected layers, each of them with a variable number of artificial neurons. A single artificial neuron with $I$ inputs ($\mathbf{x}$) calculates the output as

$$\hat{y} \equiv \sigma(\mathbf{w}^T \cdot \mathbf{x} + b)$$

that is the weighted sum of the inputs $\mathbf{x} \in \mathbb{R}^I$ with weights $\mathbf{w} \in \mathbb{R}^I$, plus a bias term $b \in \mathbb{R}$, followed by a nonlinear activation function $\sigma: \mathbb{R} \rightarrow \mathbb{R}$. The most common activation functions $\sigma(\cdot)$ are: The sigmoid $\sigma(x) \equiv \frac{1}{1+e^{-x}}$, the tanh $\sigma(x) \equiv \frac{e^x-e^{-x}}{e^x+e^{-x}}$, and the ReLU $\sigma(x) \equiv \max(0, x)$.
Thus, the distribution \( D \)imisation problem in Eq. (10) pertains to the supervised learning of vast majority of ANNs. In general, the distribution \( D \) depends on the specific layer. More concisely, the MLP is trained on all the samples of the training set \( S_{\text{tr}} \). Let us now introduce the supervised learning process. Ideally, in the training process we would like to find the parameters \( \theta^* \), \( \xi^* \) that minimises the empirical risk \( L_{S_{\text{tr}}} (\theta^*, \xi^*) \) and by the fixed hyperparameters \( \xi \) defining the number, the dimension, and the activation functions of the MLP layers.

### Supervised Training

Let us now introduce the supervised learning process. For the sake of clarity, we just refer to the training of the MLP; however, the same notions can be applied in general to the supervised learning of vast majority of ANNs.

Eq. (10) behaves like a generic function approximator \( \sigma \). Ideally, in the training process we would like to find the parameters

\[
\theta^* = \arg \min_\theta L_D (\theta, \xi)
\]

that minimise the theoretical risk function

\[
L_D (\theta, \xi) \equiv \mathbb{E}_{(x,y) \sim D} [\ell (f(x; \theta, \xi), y)],
\]

i.e., the expected value of \( \ell \) for \( (x,y) \) sampled from the distribution \( D \) that generates the data set \( \{x,y\} \). In Eq. (12), \( \ell : \mathbb{R}^{O \times O} \rightarrow \mathbb{R}^+ \) denotes the loss function (usually taken as a differentiable function, apart removable discontinuities) that measures the distance between the prediction \( \hat{y} \) and the desired output \( y \). In general, the distribution \( D \) is unknown; thus, the minimisation problem in Eq. (11) cannot be neither calculated nor solved. Indeed, one can dispose of a finite set \( S = \{ (x,y) \}_1^{n} \) of samples, to train, validate and test the ML-model. By considering the partition \( \{S_{\text{tr}}, S_{\text{va}}, S_{\text{te}}\} \) of \( S \), the theoretical risk function is approximated by the empirical risk function

\[
L_{S_{\text{tr}}} (\theta, \xi) = \frac{1}{|S_{\text{tr}}|} \sum_{(x,y) \in S_{\text{tr}}} \ell (f(x; \theta, \xi), y)
\]

that is the arithmetic mean of the loss function \( \ell \) evaluated on all the samples of the training set \( S_{\text{tr}} \). By minimising the empirical risk function \( L_{S_{\text{tr}}} (\theta, \xi) \) with respect to \( \theta \), the MLP is trained and \( \theta^* \) is obtained. Then, the validation set \( S_{\text{va}} \) is used to compute the empirical risk \( L_{S_{\text{va}}} (\theta^*, \xi) \) that takes as input the optimal parameters attained by the minimisation of \( L_{S_{\text{tr}}} \) (training stage). This procedure allows to check if the ML-model works also for unseen data. Notice that the minimisation of the training risk function \( L_{S_{\text{tr}}} (\theta, \xi) \) with respect to \( \theta \) is performed step-by-step over time. After each step (also called epoch), the validation risk \( L_{S_{\text{va}}} (\theta^*, \xi) \) is evaluated, and the minimisation procedure is stopped when the time-derivative of \( L_{S_{\text{va}}} (\theta^*, \xi) \) becomes positive for several epochs, thus showing overfitting \( \mathbb{S} \). In case such time-derivative remains negative or constant over time, the procedure is ended after a predefined number of epochs. The validation set \( S_{\text{va}} \) can be also used to explore other configurations \( \xi \) of the ML-model: this process is called hyperparameters optimization. In particular, after completing the training procedure using two different set of hyperparameters \( \xi \) and \( \xi' \), we obtain two minima \( \theta^* \) and \( \theta'^* \), and then compare \( L_{S_{\text{va}}} (\theta^*, \xi) \) with \( L_{S_{\text{va}}} (\theta'^*, \xi) \) to also choose the best hyperparameter. Finally, we use the test set \( S_{\text{te}} \) to calculate a significant metric (in our case, the classification accuracy) and report the results.

Regarding the hyperparameters optimization, it can be performed in different ways. The most basic technique is called grid search whereby the training and validation are carried out on a specific set of hyperparameters configurations. The random grid search considers configurations where each hyperparameter is randomly chosen within an a-priori fixed range of values. It has been proved to be more efficient than standard grid search \( \mathbb{S} \). A more sophisticated class of optimization methods is the Bayesian optimization \( \mathbb{S} \) that updates, after the training of each hyperparameters configuration, a Bayesian model of the validation error. The best hyperparameters configuration is thus chosen as the one allowing for the lower guess validation error.

### Minimisation algorithms

The most used optimisation algorithm to minimise Eq. (13) is the SGD \( \mathbb{S} \) and its adaptive variants, such as Adaptive Moment Estimation (ADAM) \( \mathbb{S} \), that changes the value of the learning rate \( \eta \) (i.e., the descent step) at each iteration. After having calculated the predictions \( \hat{y} \), the loss function \( \ell (\hat{y}, y) \) is propagated backwards (backpropagation) in the ANNs and its gradient in the weight space is calculated. Overall, the optimisation process consists in iteratively updating the value of the weights \( \theta \) according to the relation

\[
\theta_i = \theta_{i-1} - \eta \nabla_\theta L_{S_{b}} (\theta_{i-1}, \xi)
\]

where \( i \) is the index for the descent step and \( S_b \subset S_{\text{tr}} \) denotes the \( b \)-th set of samples, taken from the training
set and used for the computation of the gradient. If \( S_0 = S_t \), the algorithm is called batch SGD; if \( S_0 \) contains only one element is called on-line SGD; finally, the most common approach (we use it here) is mini-batch SGD that consider \( |S_0| = B \) with \( B \) a fixed dimension. Hence, the update of \( \theta \) follows the descent direction of the gradient, with a magnitude determined by the learning rate \( \eta \).

Now, let us introduce the specific loss function \( \ell \) considered in this paper. For classification problems with two or more classes, a common choice for \( \ell \) is the categorical cross entropy, which is defined as

\[
\ell(\hat{y}, y) = -\sum_{j=1}^{O} y^{(j)} \log \hat{y}^{(j)}.
\]

This function measures the dissimilarity between two or more probability distributions. Thus, to properly use the categorical cross entropy, it is convenient to choose the desired outputs \( y \) as Kronecker delta functions centered around the indices associated with each class to be classified. The model output \( \hat{y} \), instead, is normalised so that it represents a discrete probability distribution, i.e., a vector of positive elements summing to 1. This operation is obtained by using softmax as the activation function of the last layer:

\[
\sigma^{(i)}(z) \equiv \frac{e^{z^{(i)}}}{\sum_{j=1}^{O} e^{z^{(j)}}}
\]

where \( \sigma(z) \) is the vector having as elements \( \sigma^{(i)}(z) \), with \( i = 1, \ldots, O \), and \( z \) denotes the output of the last layer before the activation function.

In the experiments, the activation functions for the hidden layers of the MLP have been chosen among the sigmoid, hyperbolic tangent and rectifier functions accordingly to the hyperparameters optimization.

### Recurrent Neural Networks

A RNN is defined by the repeated relation

\[
h_t = r(x_t, h_{t-1}; \theta, \xi)
\]

where \( t \in \{1, \ldots, \tau\} \), \( h_t \in \mathbb{R}^d \) is a \( d \)-dimensional vector with \( d \) being an hyperparameter belonging to \( \xi \) and \( h_0 = \mathbf{0} \) (vector of zeros). The recurrent relation defines \( \tau \) hidden representations \( h_t \) (to be seen as a memory) of the input sequence \( \{x_1, \ldots, x_n\} \) with \( q = 1, \ldots, n \). If the function \( r \) is implemented as an MLP that takes as input the concatenation \( x_t \oplus h_{t-1} \) (usually called “vanilla RNN”), the model suffers the so-called vanishing gradient problem such that the weights of the last layers of the RNN are updated only with respect to the more recent input data. The vanishing gradient problem occurs when the backpropagation is performed on a high number of layers, as it could happen in our case with a large value of \( \tau \) (thus meaning long input sequences). In this regard, to mitigate the vanishing gradient problem, LSTM and GRU have been introduced. These methods use learned gated mechanisms, based on current input data and previous hidden representations, to control how to update the current hidden representation \( h_t \). Specifically, if LSTM is used, Eq. \ref{equation} needs to be slightly modified as

\[
s_t = v(x_t, s_{t-1}; \theta, \xi)
\]

\[
h_t = r(x_t, s_t; h_{t-1}; \theta, \xi)
\]

where \( s_0 = \mathbf{0} \) and \( v, r \) are, as usual, nonlinear functions.

### Classification with RNNs

Now, let us explain how to use the hidden representations to calculate the prediction \( \hat{y} \) in output from the ML-model. The common approach to calculate the prediction in classification problems is to use the RNN as an encoder of the sequence and to scale the dimension of the last hidden representation \( \text{h}_t[L] \oplus \text{h}_1[L] \) (in the more general case of bidirectional models) to the one of the output vector. This scaling can be done through a fully connected layer, or, more in general, by means of
Figure 3. Diagram of a bidirectional multi-layer RNN where the nonlinear function \( r \) is defined in Eqs. (19) and (20), \( a \) can be defined either with Eq. (26) or Eq. (27), \( f \) is provided by Eq. (22), and \( \oplus \) denotes concatenation. The input sequence \( x_t \), with \( t = 1, \ldots, \tau \), is processed sequentially in both directions by the function \( r \) that is parametrised by the shared sets of weights \( \theta_r[1] \) and \( \tilde{\theta}_r[1] \) for the forward and backward directions, respectively. The hidden representations \( h_t[1] \) and \( \tilde{h}_t[1] \), in turn, are processed by the subsequent layers, parametrised by a different sets of weights, so as to obtain the final hidden representations \( h_t[L] \) and \( \tilde{h}_t[L] \). Finally, \( a \) performs the aggregation of the last hidden representations adopting the attention mechanism (26) or the max pooling in Eq. (27). The classification is performed by the function \( f \) that is parametrised by \( \theta_f \). The simpler form of aggregation in Eq. (21) is not depicted in the figure.

Then, we can use SGD to minimise an empirical risk function similar to Eq. (13) of MLPs.

It is possible to consider different forms of aggregation \( a \) for the hidden representations \( h_t[L] \), with \( t = 1, \ldots, \tau \), instead of using only the last hidden representation as in Eq. (21). In this regard, attention mechanisms, also in hierarchical forms, perform a weighted average of the \( h_t[L] \) where the weights are learned together with the ML-model. In detail, Eq. (21) becomes:

\[
a = h_t[L] \oplus \tilde{h}_t[L]
\]

and

\[
\dot{y} = f(a; \theta_f, \xi).
\]

where \( \langle \cdot, \cdot \rangle \) denotes the dot product and \( c \) is a learned vector that is randomly initialised and jointly learned during the training process as in Refs. [69–72]. Another form of aggregation \( a \) is the max pooling aggregation, whereby each element \( a^{(j)} \) of \( a \) just refers to a single value of \( t \).
In this case, Eq. (21) equals to
\[ a^{(j)} = \max_t u^{(j)}_t. \] (27)
where the expression of \( u_t \) is provided by Eq. (23). In this way, each element \( u_t^{(j)} \) of the hidden representations (for \( t = 1, \ldots, \tau \)) learns to detect specific features of the input data within all the interval \([1, \tau]\).

Finally, another approach, which we do not use here, is to consider the RNN as a transducer that produces an output sequence \( \hat{y} \) for \( t = 1, \ldots, \tau \) (generally \( \tau \neq \tau \)) in correspondence of the input sequence \( x_t \) with \( t = 1, \ldots, \tau \) [43, 64, 60].

### Implementation of the machine learning algorithms

All the ML-models are realized in PyTorch and have been trained on the six different data sets using a DELL\textsuperscript{®} Precision Tower workstation with one NVIDIA\textsuperscript{®} TITAN RTX\textsuperscript{®} GPU with 10 Gb of memory, 88 cores Intel\textsuperscript{®} Xeon\textsuperscript{®} CPU E5-2699 v4 at 2.20GHz and 94 Gb of RAM.

We train the ANN models in mini-batches of dimension 16 by means of the SGD using ADAM [58] and learning rate \( \eta = 10^{-3} \). We optimize the hyperparameters \( \xi \) with ASHA [91] as scheduler and Hyperopt [92, 93] (Hyperopt belongs to the family of Bayesian optimization algorithms) as search algorithm in the framework Ray Tune [94]. For the MLP models, the hyperparameters optimization defines: (i) the activation functions to be used, (ii) the number of layers, and (iii) their dimension, within the following search space: \( \sigma \in \{\text{relu, sigmoid, tanh}\} \), \( L \in \{2, 3, 4, 5, 6\} \) and \( \dim(h[1]) \equiv \cdots \equiv \dim(h[L]) \in \{d \in \mathbb{N} | 1 \leq d \leq 512\} \). Instead, for the RNN models the search space is \( L \in \{1, 2, 3, 4\} \) for the number of recurrent layers (\( L \in \{1, 2, 3, 4, 5, 6\} \) for the NM task with \( t_{15} = 0.1 \)), and \( \dim(h[1]) \equiv \cdots \equiv \dim(h[L]) \equiv \dim(\bar{h}[L]) \in \{d \in \mathbb{N} | 1 \leq d \leq 512\} \) for the layers dimension. Regarding the ML-models \( \text{m-biGRU-att} \) and \( \text{m-biLSTM-att} \), the search space includes also the dimension of the attention layer as in Eqs. (21) and (22), i.e., \( \dim(c) \equiv \dim(v_1) \equiv \cdots \equiv \dim(v_r) \in \{d \in \mathbb{N} | 1 \leq d \leq 512\} \). In the hyperparameters optimization of all the MLP and the RNN models, we have also used regularization methods as weight decay [95] and dropout [96]. They are able to mitigate overfitting; in particular, the former adds a penalty (chosen among \( \{0, 10^{-3}, 10^{-5}\} \)) to the risk function \( L_2(\theta, \xi) \) with the aim to discourage large weights. Instead, using dropout, the outputs of the artificial neurons during the training are forced to zero with a probability among \( \{0, 0.2, 0.5\} \).

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