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# On the Identifiability of Transform Learning for Non-negative Matrix Factorization 

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#### Abstract

Non-negative matrix factorization with transform learning (TL-NMF) aims at estimating a short-time orthogonal transform that projects temporal data into a domain that is more amenable to NMF than off-the-shelf time-frequency transforms. In this work, we study the identifiability of TL-NMF under the Gaussian composite model. We prove that one can uniquely identify row-spaces of the orthogonal transform by optimizing the likelihood function of the model. This result is illustrated on a toy source separation problem which demonstrates the ability of TL-NMF to learn a suitable orthogonal basis.


Index Terms-NMF, transform learning, identifiability, source separation, joint diagonalization.

## I. Introduction

SPECTRAL unmixing by non-negative matrix factorization (NMF) is a standard approach to signal decomposition. It proceeds by transforming the signal into a domain where NMF is applied. For one-dimensional audio signals, it is customary to use a short-time frequency transform, such as the short-time Fourier or discrete cosine transforms. Such transforms first apply a short-time window to divide the signal into shorter segments of equal length and then compute the orthogonal Fourier or cosine transform separately on each segment. However, due to the time-frequency resolution trade-off, there is an increasing interest in replacing the frequency transform by a more discriminative orthogonal transform. Hence, transformlearning NMF (TL-NMF) was proposed in [1] with the goal of learning an adaptive short-time orthogonal transform. The problem is then to estimate the orthogonal transform given a fixed short-time window. TL-NMF also appears to be a special case of independent low-rank tensor analysis (ILRTA) [2] which offers a general framework for short-time modeling of temporal sequences under composite covariance models.

In this work, we study the identifiability of the orthogonal transform in TL-NMF under the Gaussian composite model (GCM) [3]. Under this probabilistic model, we first derive a negative log-likelihood objective for TL-NMF in Section 2.1. It is a variation of the Itakura-Saito divergence that is commonly used in NMF (IS-NMF) [3] and in particular in [1]. In Section 2.2 we establish conditions for identifiability of rowspaces of the orthogonal transform in TL-NMF (i.e., linear subspaces generated by subsets of its rows). The conditions generalize some results about the identifiability of jointdiagonalization proposed in [4]. Finally, Section 3 illustrates our identifiability result on a toy audio decomposition problem.

[^0]Notations: For a matrix $\mathbf{X} \in \mathbb{R}^{M \times N}$, we denote by $\mathbf{x}_{n}, \underline{\mathbf{x}}_{m}$ and $[\mathbf{X}]_{m n}\left(\right.$ or $\left.x_{m n}\right)$ its $n$-th column, $m$-th row, and $(m, n)$-th element respectively. For a vector signal $\mathbf{y} \in \mathbb{R}^{T}$, $\mathbf{y}(t)$ denotes its $t$-th element and $\mathbf{y}^{\boldsymbol{\top}}$ denotes its transpose. We write $\operatorname{Diag}(\mathbf{y})$ for the diagonal matrix formed out of the vector $\mathbf{y}$.

## II. Transform Learning NMF

Given a temporal signal $\mathbf{y} \in \mathbb{R}^{T}$, TL-NMF aims at finding a short-time orthogonal transform $\phi: \mathbb{R}^{T} \rightarrow \mathbb{R}^{M \times N}$ such that the element-wise squared magnitude $|\phi(\mathbf{y})|^{\circ 2}$ of $\phi(\mathbf{y})$ can be well approximated by a low-rank or sparse NMF, i.e.,

$$
\begin{equation*}
|\phi(\mathbf{y})|^{\circ 2} \approx \mathbf{W H} \tag{1}
\end{equation*}
$$

Denoting by $\mathbf{Y} \in \mathbb{R}^{M \times N}$ the matrix that contains $N$ adjacent and half-overlapping short-time frames of $\mathbf{y}$ with length $M$, we may write $\phi(\mathbf{y})=\boldsymbol{\Phi} \mathbf{Y}$. In this paper, we study how to identify $\boldsymbol{\Phi}$ given observations of $\mathbf{y}$ which satisfy (1). The problem of how to identify the $\mathbf{W}$ and $\mathbf{H}$ belongs to the identifiability of NMF [5]. In the next section, we derive a new objective (slightly different to [1]) for TL-NMF to learn $\boldsymbol{\Phi}$, $\mathbf{W}$, and $\mathbf{H}$. Throughout the paper, we consider $\boldsymbol{\Phi} \in \mathbb{R}^{M \times M}$ to be a real orthogonal matrix. Moreover $\mathbf{W} \in \mathbb{R}_{+}^{M \times K}$ and $\mathbf{H} \in \mathbb{R}_{+}^{K \times N}$ are non-negative matrices.

## A. Gaussian Composite Model

The Gaussian composite model (GCM) is used to characterize sound signals with composite structure [3]. The shorttime Fourier coefficients of a signal is modeled as a sum of independent Gaussian random variables with an NMF structure on their variances. We consider this model to characterize the distribution of $\phi(\mathbf{y})$. In short, under the GCM, we have

$$
\begin{equation*}
[\mathbf{\Phi} \mathbf{Y}]_{m n} \sim \mathcal{N}\left(0,[\mathbf{W H}]_{m n}\right) \tag{2}
\end{equation*}
$$

where the variance is $[\mathbf{W H}]_{m n}>0$. Moreover, conditioned on $\mathbf{\Phi}, \mathbf{W}$, and $\mathbf{H},[\mathbf{\Phi} \mathbf{Y}]_{m n}$ is assumed to be independent of $[\boldsymbol{\Phi} \mathbf{Y}]_{m^{\prime} n^{\prime}}$ for any $(m, n) \neq\left(m^{\prime}, n^{\prime}\right)$. As such, the negative log-likelihood function is given, up to a constant, by

$$
\begin{align*}
& -\log p(\mathbf{Y} \mid \mathbf{\Phi}, \mathbf{W}, \mathbf{H})=-\sum_{m n} \log \mathcal{N}\left([\mathbf{\Phi} \mathbf{Y}]_{m n}^{2} \mid 0,[\mathbf{W H}]_{m n}\right) \\
& \quad \stackrel{\text { c }}{=} \frac{1}{2} \sum_{m=1}^{M} \sum_{n=1}^{N}\left(\frac{[\mathbf{\Phi} \mathbf{Y}]_{m n}^{2}}{[\mathbf{W H}]_{m n}}+\log \left([\mathbf{W H}]_{m n}\right)\right) \tag{3}
\end{align*}
$$

When $\Phi$ is fixed (e.g., to the discrete cosine transform, DCT), minimizing (3) with respect to (w.r.t) $\mathbf{W}$ and $\mathbf{H}$ is equivalent
to IS-NMF [3]. To study identifiability, we will consider the expected (w.r.t Y) negative log-likelihood objective:

$$
\begin{aligned}
c(\mathbf{\Phi}, \mathbf{W}, \mathbf{H}) & \stackrel{\text { def }}{=} \mathbb{E}(-\log p(\mathbf{Y} \mid \mathbf{\Phi}, \mathbf{W}, \mathbf{H})) \\
& \stackrel{c}{=} \frac{1}{2} \sum_{m=1}^{M} \sum_{n=1}^{N}\left(\frac{\mathbb{E}\left([\mathbf{\Phi} \mathbf{Y}]_{m n}^{2}\right)}{[\mathbf{W H}]_{m n}}+\log \left([\mathbf{W H}]_{m n}\right)\right)
\end{aligned}
$$

## B. Identifiability of $\boldsymbol{\Phi}$

Let us assume that there exists an orthogonal transform $\bar{\Phi}$ such that $\overline{\mathbf{\Phi}} \mathbf{Y}$ follows the GCM with parameters $\overline{\mathbf{W}} \in \mathbb{R}_{+}^{M \times \bar{K}}$ and $\overline{\mathbf{H}} \in \mathbb{R}_{+}^{\bar{K} \times N}$. Here $\bar{K}$ is the ground-truth rank which is not known a priori. Identifiability is about whether the minimal solution of $c$ as of (4) in $\boldsymbol{\Phi}$ corresponds to $\overline{\boldsymbol{\Phi}}$. In the context of TL-NMF, we are interested in identifying row-spaces of $\overline{\mathbf{\Phi}}$ such that signal components of different nature can be transformed into different (orthogonal) rowspaces of $\overline{\boldsymbol{\Phi}}$. This principle can improve source separation performance as we shall illustrate in Section 3. More formally, we study identifiability in the sense of Definition 1.

Definition 1. Let $\operatorname{span}_{\bar{B}}(\bar{\Phi})$ (resp., span ${ }_{B}(\boldsymbol{\Phi})$ ) denote the row-space of $\overline{\boldsymbol{\Phi}}$ (resp., $\boldsymbol{\Phi}$ ) spanned by its rows indexed by $\bar{B}$ (resp., B). We say that the rows of $\overline{\mathbf{\Phi}}$ indexed by $\bar{B} \subset\{1, \ldots, M\}$ are identifiable if any global minimizer of $c$ is attained for a $\mathbf{\Phi}$ such that there exists $B \subset\{1, \ldots, M\}$ for which $\operatorname{span}_{\bar{B}}(\overline{\mathbf{\Phi}})=\boldsymbol{\operatorname { s p a n }}_{B}(\boldsymbol{\Phi})$.

We first give sufficient conditions for which $(\overline{\mathbf{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}})$ is a global minimum of $c$.

Lemma 1. Let $\mathbf{Y}$ be a random matrix such that, for each column $\mathbf{y}_{n}$ of $\mathbf{Y}$,

$$
\begin{align*}
\mathbb{E}\left(\mathbf{y}_{n}\right) & =\mathbf{0}  \tag{5}\\
\overline{\mathbf{\Phi}} \mathbb{E}\left(\mathbf{y}_{n} \mathbf{y}_{n}^{\top}\right) \overline{\boldsymbol{\Phi}}^{\top} & =\operatorname{Diag}\left(\overline{\mathbf{v}}_{n}\right) \tag{6}
\end{align*}
$$

for an orthogonal matrix $\overline{\mathbf{\Phi}} \in \mathbb{R}^{M \times M}$ and a vector $\overline{\mathbf{v}}_{n}=$ $\sum_{k=1}^{\bar{K}} \overline{\mathbf{w}}_{k} \bar{h}_{k n}$ formed by the non-negative matrices $\overline{\mathbf{W}} \in$ $\mathbb{R}_{+}^{M \times \bar{K}}$ and $\overline{\mathbf{H}} \in \mathbb{R}_{+}^{\bar{K} \times N}$. Then, for $K \geq \bar{K}$, the triplet $(\overline{\boldsymbol{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}})$ is a global minimizer of $c$ :

$$
\begin{equation*}
c(\boldsymbol{\Phi}, \mathbf{W}, \mathbf{H}) \geq c(\overline{\boldsymbol{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}}), \forall(\boldsymbol{\Phi}, \mathbf{W}, \mathbf{H}) \tag{7}
\end{equation*}
$$

The proof of Lemma 1 is given in Appendix A. In the GCM, each column $\mathbf{y}_{n}$ of $\mathbf{Y}$ follows $\overline{\boldsymbol{\Phi}} \mathbf{y}_{n} \sim \mathcal{N}\left(\mathbf{0}, \operatorname{Diag}\left(\overline{\mathbf{v}}_{n}\right)\right)$. Therefore the conditions (5)-(6) of Lemma 1 are always satisfied. Moreover, it is noteworthy to mention that Lemma 1 applies to a much broader class of signals than those following the GCM. For instance, when $\mathbf{Y}$ is formed by half-overlapping shorttime windows, there does not exist a $\overline{\mathbf{\Phi}}$ for which $[\overline{\mathbf{\Phi}} \mathbf{Y}]_{m n}$ is independent of $[\overline{\mathbf{\Phi}} \mathbf{Y}]_{m^{\prime} n^{\prime}}$ for adjacent $n$ and $n^{\prime}$. Hence the independence assumption from the GCM under which the objective $c$ has been derived is not satisfied. However, the conditions in Lemma 1 can still be satisfied, even if $\mathbf{Y}$ is non-Gaussian. Such examples shall be given in Section 3.

The condition (6) writes as a joint diagonalization of all the covariance matrices $\mathbb{E}\left(\mathbf{y}_{n} \mathbf{y}_{n}^{\top}\right)$, by a common basis $\overline{\boldsymbol{\Phi}}$. The identification of the rows of $\overline{\boldsymbol{\Phi}}$ is a well-studied problem in the literature of joint diagonalization [4]. However, the results are mostly about when $\boldsymbol{\Phi}$ equals to $\bar{\Phi}$ up to a signed permutation.

Below we go one step further by studying the identification of row-spaces of $\bar{\Phi}$.

Theorem 1. Let us partition the rows of $\mathbf{W H}$ into $Q$ sets $\left\{\bar{B}_{q}\right\}_{q=1}^{Q}$ with the smallest possible $Q$, such that $\forall\left(m, m^{\prime}\right) \in$ $\left(\bar{B}_{q}\right)^{2}$ and $\forall n \in\{1, \ldots, N\}, \quad[\overline{\mathbf{W}} \overline{\mathbf{H}}]_{m, n}=[\overline{\mathbf{W}} \overline{\mathbf{H}}]_{m^{\prime}, n}$ (partition into subsets of equal rows, if any). Then, under the conditions of Lemma 1 (i.e. (5), (6), and $K \geq \bar{K}$ ), $\forall q \in\{1, \ldots, Q\}$ the rows of $\bar{\Phi}$ indexed by $\bar{B}_{q}$ are identifiable.

The proof of Theorem 1 is given in Appendix B. Theorem 1 coincides with Theorem 2 in [4] when $Q=M$. Our proof technique is different from [4] as we analyze the global minimizers of $c$, although they both rely on the orthogonality assumption of $\boldsymbol{\Phi}$.

Numerical illustration: Let us provide a numerical example to minimize $c$ when $\mathbf{Y}$ follows the GCM with the ground-truth parameters $\overline{\boldsymbol{\Phi}}, \overline{\mathbf{W}}$, and $\overline{\mathbf{H}}$. In such case, the expectation in (4) is given by (see Appendix A)

$$
\begin{equation*}
\mathbb{E}\left(|\boldsymbol{\Phi} \mathbf{Y}|^{\circ 2}\right)=\left(\boldsymbol{\Phi} \overline{\boldsymbol{\Phi}}^{\boldsymbol{\top}}\right)^{\circ 2} \overline{\mathbf{W}} \overline{\mathbf{H}} \tag{8}
\end{equation*}
$$

Let us consider the case where $M=4, N=3, \bar{K}=K=2$, $\bar{\Phi}$ is a type-III DCT and

$$
\overline{\mathbf{W}} \overline{\mathbf{H}}=\left(\begin{array}{ccc}
1 & 0.1 & 1.1  \tag{9}\\
1 & 0.1 & 1.1 \\
0.1 & 1 & 1.1 \\
0.1 & 1 & 1.1
\end{array}\right)
$$

Given the structure of $\overline{\mathbf{W}} \overline{\mathbf{H}}$ and Theorem 1, there are $Q=2$ identifiable row-spaces of $\overline{\boldsymbol{\Phi}}$. The first (resp., second) one is the linear span of the first and second rows (resp., the third and fourth rows) of $\overline{\boldsymbol{\Phi}}$ (i.e., $\bar{B}_{1}=\{1,2\}$ and $\bar{B}_{2}=\{3,4\}$ ). To minimize $c$, we use a slight modification (that reflects the GCM assumption instead of pure IS divergence) of the blockcoordinate descent algorithm described in [6] (projected quasiNewton step for $\boldsymbol{\Phi}$, majorization-minimization for $\mathbf{W}$ and $\mathbf{H}$ ). The algorithm starts from a random initialization and stops when sufficient numerical precision is reached. The estimated basis $\boldsymbol{\Phi}$ satisfies

$$
\boldsymbol{\Phi} \overline{\boldsymbol{\Phi}}^{\boldsymbol{\top}}=\left(\begin{array}{cccc}
0.703 & 0.711 & 0.000 & 0.000 \\
0.000 & 0.000 & 0.955 & -0.297 \\
0.711 & -0.703 & 0.000 & 0.000 \\
0.000 & 0.000 & -0.297 & -0.955
\end{array}\right)
$$

This shows that the first and second rows (resp., the third and forth row) of $\bar{\Phi}$ are accurately identified with the first and third row (resp., the second and forth row) of $\boldsymbol{\Phi}$. We remark that the condition (6) in Lemma 1 implies that existing joint-diagonalization algorithms [7], [8], [9], [10] could also be applicable to identify $\overline{\boldsymbol{\Phi}}$ directly, prior to NMF.

## III. Unsupervised Source Separation

We evaluate the performance of TL-NMF under the GCM in an unsupervised single-channel source separation problem. We study whether TL-NMF can learn a better $\boldsymbol{\Phi}$ than fixed DCT to improve source separation performance. The code to reproduce the numerical experiments reported in this section is available online. ${ }^{1}$

[^1]a) Experimental Setting: We consider the signal $\mathbf{y}(t)=$ $\sum_{i=1}^{I} \mathbf{y}^{(i)}(t)$ where $\left(\mathbf{y}^{(i)}\right)_{i=1}^{I}$ are random sources defined as the sum of cosine functions shifted by a random phase to mimic musical notes with different timbres. More precisely, we set $\mathbf{y}^{(i)}(t)=\sum_{r=1}^{R} \beta_{i, r} \cos \left(r\left(\alpha_{i} t+\theta_{i}\right)\right) \mathbf{g}_{i}(t)$ where $\theta_{i} \in$ $[0,2 \pi)$ is an uniform random phase, $\alpha_{i}>0$ a base frequency, $r$ an integer exponent, $\beta_{i, r}>0$ a timbre coefficient, and $\mathbf{g}_{i} \in \mathbb{R}^{T}$ a positive envelop that varies slowly over $t \in\{1, \ldots, T\}$. Moreover, we assume that $\theta_{i}$ is independent of $\theta_{i^{\prime}}$, for all $i^{\prime} \neq i$. Our goal is then to assess the ability of TL-NMF to separate those $I$ harmonic notes from the signal $\mathbf{y}$.

Let us analyse to which extent the assumptions of Lemma 1 are satisfied by the considered signal $\mathbf{y}$. Given a short-time window $\mathbf{a} \in \mathbb{R}^{M}$, the short-time matrix $\mathbf{Y}^{(i)}$ for the $i$-th source verifies $\mathbf{y}_{n}^{(i)}(m)=\mathbf{y}^{(i)}\left(m+\frac{M n}{2}\right) \mathbf{a}(m)$ (using zero-padding at signal boundaries). From the independence of the uniform random phases $\left\{\theta_{i}\right\}_{i=1}^{I}$, we get that $\mathbb{E}\left(\mathbf{y}_{n}\right)=\sum_{i} \mathbb{E}\left(\mathbf{y}_{n}^{(i)}\right)=\mathbf{0}$ for all $n \in\{1, \ldots, N\}$. Similarly, $\mathbb{E}\left(\mathbf{y}_{n} \mathbf{y}_{n}^{\boldsymbol{\top}}\right)=\sum_{i} \mathbb{E}\left(\mathbf{y}_{n}^{(i)} \mathbf{y}_{n}^{(i) \mathbf{T}}\right)$ and its $\left(m, m^{\prime}\right)$-th element reads

$$
\begin{equation*}
\sum_{i=1}^{I} \frac{s_{m, m^{\prime}, n}^{(i)}}{2} \sum_{r=1}^{R} \beta_{i, r}^{2} \cos \left(r \alpha_{i}\left(m-m^{\prime}\right)\right) \tag{10}
\end{equation*}
$$

where $s_{m, m^{\prime}, n}^{(i)}=\mathbf{a}(m) \mathbf{a}\left(m^{\prime}\right) \mathbf{g}_{i}\left(m+\frac{M n}{2}\right) \mathbf{g}_{i}\left(m^{\prime}+\frac{M n}{2}\right)$. Hence, the question is whether there exist $\overline{\mathbf{\Phi}}, \overline{\mathbf{W}}$, and $\overline{\mathbf{H}}$, such that $\mathbb{E}\left(\mathbf{y}_{n} \mathbf{y}_{n}^{T}\right)$ satisfies the joint-diagonalization condition (6). A general analysis from (10) is challenging. Instead, let us provide an intuition through the special case where $s_{m, m^{\prime}, n}^{(i)}=s_{n}^{(i)}$ (i.e., independent from $m$ and $m^{\prime}$ ) and, $\forall i$, $r \alpha_{i}=r \alpha=\pi p / M$ for $p \in \mathbb{N} \cap(0, M)$. Then, one easily gets that $\mathbb{E}\left(\mathbf{y}_{n} \mathbf{y}_{n}^{\boldsymbol{\top}}\right)=\overline{\boldsymbol{\Phi}}^{\boldsymbol{\top}} \operatorname{Diag}\left(\overline{\mathbf{v}}_{n}\right) \overline{\boldsymbol{\Phi}}$ for

$$
\begin{align*}
& \overline{\mathbf{v}}_{n}=\sum_{i=1}^{I} \frac{M s_{n}^{(i)}}{4}\left(\beta_{i, 1}^{2}, \beta_{i, 1}^{2}, \cdots, \beta_{i, R}^{2}, \beta_{i, R}^{2}, 0 \cdots 0\right)^{\top} \in \mathbb{R}_{+}^{M} \\
& \overline{\mathbf{\Phi}}=\left(\mathbf{c}_{1}, \mathbf{s}_{1}, \cdots, \mathbf{c}_{R}, \mathbf{s}_{R}, * \cdots *\right)^{\top} \in \mathbb{R}^{M \times M} \tag{11}
\end{align*}
$$

where $\mathbf{c}_{r}(m)=\sqrt{2 / M} \cos \left(r \alpha m+\varphi_{r}\right)$ and $\mathbf{s}_{r}(m)=$ $\sqrt{2 / M} \sin \left(r \alpha m+\varphi_{r}\right)$ for some phase shift $\varphi_{r} \in[0,2 \pi)$. The last $M-2 R$ atoms are arbitrary as long as $\overline{\boldsymbol{\Phi}}$ is orthogonal. Hence, the assumptions of Lemma 1 are satisfied with $\bar{K}=I$. ${ }^{2}$

In this work, we consider a scenario with $I=2$ sources that are generated according to the aforementioned model with $R=2, \beta_{1, r}=\beta_{2, r}=0.5^{r}$, as well as the base periods $\frac{2 \pi}{\alpha_{1}}=\frac{40}{\pi} \mathrm{~ms}$ and $\frac{2 \pi}{\alpha_{2}}=\frac{25}{\pi} \mathrm{~ms}$, respectively. The signal duration is 3 s with sampling rate 5000 Hz , which is large enough to avoid aliasing. One realization of such a signal is depicted in Figure 1, where one can also see the shape of the used envelopes $\mathbf{g}_{i}$. To deploy our TL-NMF method, we set a to be a Tukey window [11] with a cosine fraction parameter set to 0.1. The duration of the window is fixed to $40 \mathrm{~ms}(M=200)$. As such, each segment covers around 3 or 5 periods of each source. Although this configuration does not follow the (ideal) assumptions we made to derive (11), the covariance matrix $\mathbb{E}\left(\mathbf{y}_{n} \mathbf{y}_{n}^{\boldsymbol{\top}}\right)$ is still close to be circulant Toeplitz. Indeed, given the very slow variations of the envelopes $\mathbf{g}_{i}$ with respect to the short duration of the Tukey window $\mathbf{a}(m) \approx 1$, we have $s_{m, m^{\prime}, n}^{(i)} \approx \mathbf{g}_{i}\left(\frac{M n}{2}\right)^{2}$. Hence, we expect that TL-NMF learns

[^2]

Fig. 1: One random realization $\mathbf{y}$ of a sum of two harmonic processes. Source $\mathbf{y}_{1}(t)$ (resp., $\mathbf{y}_{2}(t)$ ) is nonzero for $t \in[0,1] \cup[2,3]$ (resp., $t \in[1,3]$ ).
a transform $\boldsymbol{\Phi}$ containing atoms alike $\cos \left(r \alpha_{i} m+\varphi_{r, i}\right)$ and $\sin \left(r \alpha_{i} m+\varphi_{r, i}\right)$. As $s_{m, m^{\prime}, n}^{(i)}$ depends on a, the choice of the window will have an impact on the learnt basis $\boldsymbol{\Phi}$. This phenomenon has been previously observed in [1], [2], [12].
b) Estimation: Unlike the ideal case where $c$ can be computed from (4) and (8), we generate 10 i.i.d realizations of the mixed signal $y$ to construct an empirical estimate of $\mathbb{E}\left(|\Phi \mathbf{Y}|^{\circ 2}\right)$, and thus of $c$. This allows us to apply the numerical method deployed in II.B to optimize $\boldsymbol{\Phi}, \mathbf{W}$, and H. As a baseline, we consider IS-NMF under a fixed typeIII DCT basis. It amounts to minimize the estimated $c$ with respect to $\mathbf{W}$ and $\mathbf{H}$ under the fixed $\boldsymbol{\Phi}$.
c) Results: Figure 2 displays the NMF of $\mathbb{E}\left(|\Phi \mathbf{Y}|^{\circ 2}\right)$ using $K=2$. We see that the columns $\mathbf{w}_{k}$ of $\mathbf{W}$ are better separated under the $\boldsymbol{\Phi}$ from TL-NMF than under the DCT. Each row $\underline{\mathbf{h}}_{k}$ of $\mathbf{H}$ varies also smoother over $n$, following the envelope $\mathbf{g}_{i}(t)$ of each source. The eight rows of $\boldsymbol{\Phi}$ with largest energy $\sum_{n} \mathbb{E}\left([\boldsymbol{\Phi} \mathbf{Y}]_{m n}^{2}\right)$ are depicted in Figure 3. It illustrates that the TL-NMF has found atoms alike $\cos \left(r \alpha_{i} t+\varphi_{i, r}\right)$ or $\sin \left(r \alpha_{i} t+\varphi_{i, r}\right)$, as expected.

Finally, we apply standard Wiener filtering to separate each realization of $\mathbf{y}$ into $K=2$ sources. It involves computing two masks from the estimated $\mathbf{W}$ and $\mathbf{H}$ to be applied to $\Phi \mathbf{Y}$ prior to overlap-add reconstruction, see [3]. Quality of separation is evaluated with BSS_Eval 2.1 [13]. We evaluate the separation of the ten realizations of $\mathbf{y}$ and average the BSS_Eval metrics. We then repeat this experiment 5 times and report the mean values with standard deviation of the metrics in Table I. They show that the performance is significantly improved over DCT, thanks to the learning of a more meaningful basis $\mathbf{\Phi}$.

|  | source | SDR | SIR | SAR |
| :---: | :---: | :---: | :---: | :---: |
| IS-NMF | 1 | $16.00(0.04)$ | $22.64(0.11)$ | $17.10(0.07)$ |
|  | 2 | $16.13(0.07)$ | $23.53(0.17)$ | $17.02(0.05)$ |
| TL-NMF | 1 | $22.53(0.16)$ | $40.20(0.99)$ | $22.61(0.18)$ |
|  | 2 | $22.52(0.17)$ | $39.00(0.23)$ | $22.62(0.17)$ |

TABLE I: BSS_Eval mean values with standard deviation (in bracket).

## IV. CONCLUSION

In this paper, we have studied the identifiability of transform learning for NMF under the Gaussian composite model (GCM). By minimizing the expected negative log-likelihood, we prove that it is possible to uniquely identify the row-spaces of $\bar{\Phi}$. This identifiability result is further supported by a source separation example which demonstrates the ability of TL-NMF to separate musical notes with close base frequencies. We may extend our identifiability results to the more general ILRTA [2] in future work.


Fig. 2: NMF factors $\mathbf{W}$ and $\mathbf{H}$ of $\mathbb{E}\left(|\Phi \mathbf{Y}|^{\circ 2}\right)$. Top: with DCT. Bottom: with $\boldsymbol{\Phi}$ from TL-NMF. The $m$-axis is reordered with decreasing energy $\sum_{n} \mathbb{E}\left([\boldsymbol{\Phi} \mathbf{Y}]_{m n}^{2}\right)$. We zoom in the $m$-axis where $w_{k m}$ is non-negligible.


Fig. 3: The eight most significant atoms of $\boldsymbol{\Phi}$ from TL-NMF.

## APPENDIX

## A. Proof of Lemma 1

To show that $c(\mathbf{\Phi}, \mathbf{W}, \mathbf{H}) \geq c(\overline{\mathbf{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}})$, we first minimize $c$ with respect to $\mathbf{W}, \mathbf{H}$. For a fixed orthogonal $\boldsymbol{\Phi} \in \mathbb{R}^{M \times M}$, let us first show that

$$
\begin{equation*}
c^{*}(\boldsymbol{\Phi})=\min _{\mathbf{W}, \mathbf{H}} c(\mathbf{\Phi}, \mathbf{W}, \mathbf{H})=\sum_{m n}\left(1+\log \mathbb{E}\left([\boldsymbol{\Phi} \mathbf{Y}]_{m n}^{2}\right)\right) \tag{12}
\end{equation*}
$$

and that this minimal value is attained at any point $(\mathbf{W}, \mathbf{H})$ that verifies $\mathbf{W H}=\mathbb{E}\left(|\boldsymbol{\Phi Y}|^{\circ 2}\right)$. To that end, it is sufficient to show that, under conditions (5)-(6), there exists an exact NMF factorization WH of $\mathbb{E}\left(|\boldsymbol{\Phi} \mathbf{Y}|^{\circ 2}\right)$ for any $\boldsymbol{\Phi}$. Let $\mathbf{D}=\boldsymbol{\Phi} \overline{\boldsymbol{\Phi}}^{\boldsymbol{\top}}$, so that $\boldsymbol{\Phi}=\mathbf{D} \bar{\Phi}$. Denote $\overline{\mathbf{V}}=\overline{\mathbf{W}} \overline{\mathbf{H}}$. From (5) and (6) we have

$$
\begin{align*}
\mathbb{E}\left([\mathbf{\Phi} \mathbf{Y}]_{m n}^{2}\right) & =\mathbb{E}\left(\left(\sum_{m^{\prime}} d_{m m^{\prime}}[\overline{\mathbf{\Phi}} \mathbf{Y}]_{m^{\prime} n}\right)^{2}\right)  \tag{13}\\
& =\sum_{m^{\prime}} d_{m m^{\prime}}^{2} \bar{v}_{m^{\prime} n} \tag{14}
\end{align*}
$$

Hence, because all the terms in (14) are non-negative, we have derived an exact NNF factorization of $\mathbb{E}\left(|\boldsymbol{\Phi} \mathbf{Y}|^{\circ 2}\right)$ which proves (12). As $\boldsymbol{\Phi}$ is orthogonal, one can verify that the objective $c^{*}(\boldsymbol{\Phi})$ is equivalent to the one used in [8], [10].

Then, it follows that

$$
\begin{equation*}
c^{*}(\boldsymbol{\Phi})=\sum_{m n} \log \left(\left[\left(\boldsymbol{\Phi} \overline{\boldsymbol{\Phi}}^{\mathbf{\top}}\right)^{\circ 2} \overline{\mathbf{V}}\right]_{m n}\right)+N M \tag{15}
\end{equation*}
$$

It remains to prove that $\min _{\boldsymbol{\Phi}} c^{*}(\boldsymbol{\Phi})$ is attained at $\boldsymbol{\Phi}=\overline{\boldsymbol{\Phi}}$. Using the fact that $\mathbf{D}=\boldsymbol{\Phi} \bar{\Phi}^{\top}$ is orthogonal and thus that both the columns and the rows of $\mathbf{D}^{\circ 2}$ sum to one, we obtain from Jensen's inequality

$$
\begin{equation*}
\log \left(\sum_{m^{\prime}} d_{m m^{\prime}}^{2} \bar{v}_{m^{\prime} n}\right) \geq \sum_{m^{\prime}} d_{m m^{\prime}}^{2} \log \left(\bar{v}_{m^{\prime} n}\right) \tag{16}
\end{equation*}
$$

This implies the following lower bound on $c^{*}(\boldsymbol{\Phi})$

$$
\begin{align*}
c^{*}(\boldsymbol{\Phi}) & \geq \sum_{m n} \sum_{m^{\prime}} d_{m m^{\prime}}^{2} \log \left(\bar{v}_{m^{\prime} n}\right)+N M  \tag{17}\\
& =\sum_{m^{\prime} n} \log \left(\bar{v}_{m^{\prime} n}\right)+N M=c(\overline{\mathbf{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}}), \tag{18}
\end{align*}
$$

and completes the proof.

## B. Proof of Theorem 1

Let $(\boldsymbol{\Phi}, \mathbf{W}, \mathbf{H})$ be a global minimizer of $c$ such that $c(\boldsymbol{\Phi}, \mathbf{W}, \mathbf{H})=c(\overline{\boldsymbol{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}})$. Following the results in Lemma 1, we are going to show that $\forall q \in\{1, \ldots, Q\}$, there exists a partition $\left\{B_{q}\right\}_{q=1}^{Q}$ of $\{1, \ldots, M\}$ such that $\boldsymbol{\operatorname { s p a n }}_{B_{q}}(\boldsymbol{\Phi})=\boldsymbol{\operatorname { s p a n }}_{\bar{B}_{q}}(\overline{\boldsymbol{\Phi}})$. By definition of $\bar{B}_{q}$, we denote $\hat{v}_{q n}=\bar{v}_{m n}$ for $m \in \bar{B}_{q}$. Hence $\forall q^{\prime} \neq q$, there exists $n \in\{1, \ldots, N\}$ such that $\hat{v}_{q n} \neq \hat{v}_{q^{\prime} n}$. As $c(\mathbf{\Phi}, \mathbf{W}, \mathbf{H})=$ $c(\overline{\mathbf{\Phi}}, \overline{\mathbf{W}}, \overline{\mathbf{H}})$, the equality holds in (16) and we can write it as

$$
\begin{equation*}
\forall n, \quad \log \left(\sum_{q=1}^{Q} \hat{d}_{m q}^{2} \hat{v}_{q n}\right)=\sum_{q=1}^{Q} \hat{d}_{m q}^{2} \log \left(\hat{v}_{q n}\right) \tag{19}
\end{equation*}
$$

where $\hat{d}_{m q}^{2}=\sum_{m^{\prime} \in \bar{B}_{q}} d_{m m^{\prime}}^{2}$. We will show at the end of the proof that $\forall m \in\{1, \ldots, M\}$ there exists $\tau(m) \in\{1, \ldots, Q\}$ such that

$$
\hat{d}_{m q}^{2}= \begin{cases}1 & \text { if } q=\tau(m)  \tag{20}\\ 0 & \text { if } q \neq \tau(m)\end{cases}
$$

From (20), we can construct $B_{q}=\{m: \tau(m)=q, 1 \leq m \leq$ $M\}$ so that

$$
\begin{equation*}
\hat{d}_{m q}^{2}=\sum_{m^{\prime} \in \bar{B}_{q}} d_{m m^{\prime}}^{2}=1 \quad \text { if } \quad m \in B_{q} . \tag{21}
\end{equation*}
$$

By the definition of $d_{m m^{\prime}}$ and the orthogonality of $\boldsymbol{\Phi}$ and $\overline{\boldsymbol{\Phi}}$, (21) implies that each row of $\Phi$ in the set $B_{q}$ belong to $\boldsymbol{\operatorname { s p a n }}_{\bar{B}_{q}}(\overline{\boldsymbol{\Phi}})$. Therefore $\boldsymbol{\operatorname { s p a n }}_{B_{q}}(\boldsymbol{\Phi}) \subset \boldsymbol{\operatorname { s p a n }}_{\bar{B}_{q}}(\overline{\boldsymbol{\Phi}})$.

To show that $\boldsymbol{\operatorname { s p a n }}_{B_{q}}(\boldsymbol{\Phi})=\boldsymbol{\operatorname { s p a n }}_{\bar{B}_{q}}(\bar{\Phi})$, we first check that

$$
\begin{equation*}
\forall q, \quad\left|B_{q}\right|=\sum_{m=1}^{M} \hat{d}_{m q}^{2}=\sum_{m^{\prime} \in \bar{B}_{q}} \sum_{m=1}^{M} d_{m m^{\prime}}^{2}=\left|\bar{B}_{q}\right| . \tag{22}
\end{equation*}
$$

From (22), we also obtain a similar formula as (21)

$$
\begin{equation*}
\sum_{m \in B_{q}} d_{m m^{\prime}}^{2}=1 \quad \text { if } \quad m^{\prime} \in \bar{B}_{q} \tag{23}
\end{equation*}
$$

Therefore $\boldsymbol{\operatorname { s p a n }}_{\bar{B}_{q}}(\overline{\boldsymbol{\Phi}}) \subset \operatorname{span}_{B_{q}}(\boldsymbol{\Phi})$.
Finally we show that (20) is correct by absurd. Assume that (20) is not true. Hence there exist $m$ and $q^{\prime} \neq q$, such that

$$
\begin{equation*}
\hat{d}_{m q}^{2} \in(0,1) \text { and } \hat{d}_{m q^{\prime}}^{2} \in(0,1) \tag{24}
\end{equation*}
$$

Then, the fact that $\sum_{q} \hat{d}_{m q}^{2}=\sum_{m^{\prime}} d_{m m^{\prime}}^{2}=1$ together with the strict concavity of the $\log$ function imply that the equality (19) can hold only if $\hat{v}_{q n}=\hat{v}_{q^{\prime} n}, \forall n$. This contradicts the fact that $\bar{B}_{q} \cap \bar{B}_{q^{\prime}}=\emptyset$ and completes the proof.

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[^1]:    ${ }^{1}$ to be released at publication.

[^2]:    ${ }^{2}$ When $\overline{\mathbf{v}}_{n}$ contains zeros as in (11), we add a small constant $\epsilon$ to each $[\boldsymbol{\Phi} \mathbf{Y}]_{m n}$ and $[\mathbf{W H}]_{m n}$ in the objective (4). Here, we set $\epsilon=5 \times 10^{-7}$.

