# <span id="page-0-0"></span>Structural Modelling of Dynamic Networks and Identifying Maximum Likelihood

Gourieroux<sup>∗</sup> , C. and J. Jasiak†

October 5, 2024

#### Abstract

This paper examines nonlinear dynamic models where the main parameter of interest is a nonnegative matrix characterizing the network (contagion) effects. This network matrix is assumed either to have a limited number of nonzero elements (sparsity), or to admit a reduced-rank Nonnegative Matrix Factorization (NMF). We follow the latter approach and develop new probabilistic NMF inference methods. We introduce a novel Identifying Maximum Likelihood (IML) method for consistent estimation of the identified set of admissible NMF's and derive the asymptotic distribution of this random set. We also propose a maximum likelihood estimator of the parameter matrix, for a given non-negative rank, and derive its asymptotic distribution and the associated efficiency bound.

Keywords: Network, Nonnegative Matrix Factorization (NMF), Set Identification, Random Set, Heterogeneity, Identifying Maximum Likelihood (IML), Alternating Maximum Likelihood (AML).

<sup>∗</sup>University of Toronto, Toulouse School of Economics and CREST, e-mail:Christian.Gourieroux@ENSAE.fr.

<sup>†</sup>York University, e-mail: jasiakj@yorku.ca.

We thank C. Bontemps, X. D'Haultfoeuille, R. Engle, M. Henry, A. Monfort and the participants of the CREST Financial Econometrics seminar, CIREQ Econometrics Conference in Montreal 2024, African Meeting of the Econometric Society, TSE Econometrics Seminar and Panel Data Conference in Orleans for helpful comments.

## 1 Introduction

A network is commonly characterized by a non-negative matrix A of dimension  $(n, m)$ . There exists a large literature on network models and different assumptions are imposed on the non-negative matrix A. These assumptions concern either the 0 - 1 entries of adjacency matrices, or the positive elements and "reduced rank" of interaction matrices. In applications to contagion, transmission, social interaction, and spillover effects, the elements of each row of matrix  $A$  can be constrained to sum up to one. Then, each row of matrix  $A$  is interpreted as a conditional probability distribution and matrix A is called a transition, or migration matrix.

Depending on the objective of analysis, network models are either predictionoriented, or structural. The prediction-oriented models are used for image and facial recognition, text analysis, environmetrics, and genomics. In these applications, the dimensions n and m of matrix A are very large, and the primary objective of the analysis is to reduce this high dimensionality. Then, adjacency matrices with small numbers of links, i.e. entries equal to 1, are used to obtain sparse matrices A. Alternatively, interaction matrices are chosen so that the non-negative matrix factorization (NMF) of matrix A is of a small reduced rank [Berman and Plemmons (1994)].

The network analysis in econometrics is more structural and parametric [see Manski (1993) for the reflection problem, Blume et al. (2011), de Paula (2017) for a survey and the references therein]. The panel models often represent an equilibrium, such as:  $Y_t = A Y_t + D X_t + u_t$ , where  $Y_t = (y_{1,t},..., y_{n,t})$  is the vector of individual observations at time t, and  $X_t$  are the explanatory variables. At equilibrium, a simultaneity arises because variable  $Y$  appears on both sides of the equation. The dimension  $(n, n)$  of matrix A is much smaller than in image analysis, for example. Moreover, matrix A is often defined as a parametric non-negative combination  $A=\sum$ L  $_{l=1}$  $\alpha_l A_l, \ \alpha_l \geq 0, \ l = 1, \ldots, L$ , of known network matrices  $A_l$ 's in order to obtain a simple parametric model, which is a crucial assumption in this literature [see e.g. Bramoulle et al. (2009), Lee et al. (2010), Cohen-Cole et al. (2014), Blume et al.  $(2015)$ .

The aim of this paper is to fill the gap between the prediction-oriented and structural analysis by introducing the Identifying Maximum Likelihood (IML) estimation method, that provides the Maximum Likelihood (ML) estimator of the set of admissible NMFs, as well as an ML estimator of matrix A, for an extended class of dynamic probabilistic network models.

We consider a class of dynamic parametric models with non-negative interaction matrices of reduced rank. The static models used in machine learning, dynamic panel models for individual qualitative variables, and multivariate dynamic Poisson models with contagion used in epidemiology, arise as special cases in that class of models [Gourieroux and Lu (2019), (2023), Lu et al. (2023)]. The NMF of matrix A in this class of models may involve not only the principal directions, but also a latent heterogeneity distribution. In general, there exist multiple admissible factorizations of a given true network matrix  $A_0$  with or without zero entries. We describe analytically the identified set of admissible NMFs for non-negative ranks  $K = 1$  and  $K = 2$  of matrix A, and show how the identified set can be parametrized in the general case  $K \geq$ 2. Our contribution is the Identifying Maximum Likelihood (IML) approach which is introduced to estimate the identified set of NMF's and to derive the asymptotic distribution of the estimated identified set. In this context, the estimation of a NMF with the most concentrated latent heterogeneity and/or the least collinear principal directions is considered. Moreover, we propose a maximum likelihood estimator of the identified non-negative matrix A, given its non-negative rank, and derive the efficiency bound for this rank-constrained parameter matrix estimator.

Our approach differs from the methods used in the existing network literature with respect to the following: 1) Our network is examined in a nonlinear dynamic framework, and distinguishes the short and long-run effects. Our approach is applicable to static linear models existing in the literature, which arise as special cases of the considered class of models; 2) Matrix  $A$  is compatible with the existing empirical applications concerning mainly semi-aggregated level networks with matrices A without zero elements (intricate network) and with diagonal (resp. off-diagonal) elements representing the within (resp. between) segment connections; 3) In our paper, matrix A is only assumed to satisfy a factor decomposition into non-negative factorial directions<sup>[3](#page-0-0)</sup>. To solve the partial identification issue in NMF, we derive analytically the identified set and show that it can be locally parametrized. In particular, the

<sup>3</sup> i.e. without imposing a priori strong identifying assumptions.

identified set is not defined from inequality moment restrictions, as it is commonly done in the micro-econometric literature on partial identification [see Chernozhukov et al. (2007), Andrews and Soares (2010), Canay and Shaikh (2017) for a statistical approach to this partial identification]. Instead, the Identifying Maximum Likelihood (IML) method is introduced in this paper, allowing for sharp estimation of the identified set and of the non-negative rank-constrained matrix A and derivation of the asymptotic distribution of the maximum likelihood estimator of the identified set.

The paper is organized as follows. In Section 2, we describe the class of dynamic parametric models with interaction matrices of reduced rank. Section 3 discusses the identification of the NMF of matrix A. Statistical inference is developed in Section 4. Section 5 describes the relationship between the proposed method and nonparametric partial identification. An illustration is provided in Section 6. Section 7 concludes. The list of additional regularity assumptions for asymptotic statistical inference and additional estimation results are given in on-line appendices.

## 2 Parametric Model with Interaction Matrix

### 2.1 The model

We consider a set of observations  $Y_t, t = 1, \ldots, T$ , that can be scalars, vectors, or matrices. We assume that  $(Y_t)$  is a stationary Markov process and introduce a parametric model of the conditional distribution of  $Y_t$  given its lagged values, with a conditional probability density function (p.d.f.) :

$$
l(y_t|y_{t-1}) = l(y_t|y_{t-1};A),
$$
\n(2.1)

where A is an unknown non-negative matrix  $A \geq 0$  of dimension  $(n, m)$ , which has non-negative entries. We make the following assumption :

**Assumption A.1**: i) The parametric model is well-specified, with a true value  $A_0$  of matrix parameter A; ii) The process  $(Y_t)$  is strictly stationary, geometrically ergodic.

We know that a non-negative matrix  $A$  can be factorized and written as :

$$
A = BC',\tag{2.2}
$$

where B (resp. C) is of dimension  $(n, K)$  [resp.  $(m, K)$ ] and B and C are nonnegative :  $B \geq 0, C \geq 0$ . Among the multiple non-negative factorizations available, some correspond to the minimal order  $K$ , called the non-negative rank of matrix  $A$  and denoted by  $Rk_{+}(A)$ . The non-negative rank of A is always larger or equal to the rank of A.

Specifically, a non-negative matrix factorization (NMF) can be written as follows. Let  $\beta_k, k = 1, ..., K$  (resp.  $\gamma_k, k = 1, ..., K$ ) denote the columns of B (resp. C). These columns define the factorial directions. We have :

$$
B = (\beta_1, \dots, \beta_K), C = (\gamma_1, \dots, \gamma_K),
$$
\n
$$
(2.3)
$$

and then :

$$
A = \sum_{k=1}^{K} \beta_k \gamma'_k \text{ with } \beta_k \ge 0, \gamma_k \ge 0, \forall k. \tag{2.4}
$$

This is a decomposition of  $A$  into a sum of  $K$  non-negative matrices of rank 1.

In structural models, we may not only be interested to find the true value  $A_0$ , but also a true NMF :  $B_0 C'_0 = \sum$  $_{K_0}$  $k=1$  $\beta_{0,k}\gamma'_{0,k}$ , that generates  $A_0$ . We introduce the additional assumptions :

**Assumption A.2**: i) The nonnegative rank  $K_0 = Rk_+(A_0)$  is known; ii) The true matrix  $A_0$  is asymptotically identifiable, i.e. the maximization  $\max_A E_0 \log l(Y_t|Y_{t-1};A)$ with respect to the set of non-negative matrices, has a unique solution  $A = A_0$ ; iii) The vectors  $\beta_{0,k}, \gamma_{0,k}, k = 1, \ldots, K$ , have strictly positive entries.

**Remark 1**: This analysis can be easily extended to a) parametric conditional models including observed exogenous variables  $X_t$ , or b) additional parameters  $\theta$  in models of the type:  $l(y_t|y_{t-1}, x_t) = l(y_t|y_{t-1}, x_t; A_0, \theta_0)$ , and c) dynamic panel models with covariates.

### 2.2 Examples

In general, the NMF is applied without assuming a probabilistic structure. Our objective is to extend the NMF to applications, to which the principal component analysis (PCA) by Tipping, Bishop (1999) is not applicable because it does not ensure the positivity condition. As examples, we show below the models for facial recognition, epidemiology, volatility, credit risk and cyber risk analysis. These models have to account for the non-negativity of observations  $Y_t$  in practice.

#### 2.2.1 Static Model

The static models assume that observations  $Y_t, t = 1, ..., T$ , are independent and identically distributed (i.i.d.). These models are commonly used for image analysis [see Paatero, Tapper (1994) for the introduction of NMF to the applied literature]. NMF started to be extensively studied following Lee and Seung (1999) and their application to learning the parts of objects under a non-probabilistic approach. In image analysis, the observations are matrices  $Y_t = (Y_{i,j,t})$ , where  $(i, j)$  denote the coordinates of a point in picture t and the value  $Y_{i,j,t}$  provides the pixel intensities associated with coordinates  $i, j$  and picture t. The pixel intensity can be measured on either a discrete, or continuous scale. Then, the model:

$$
l(y_t|y_{t-1};A) = \Pi_{i=1}^n \Pi_{j=1}^m f(y_{i,j,t};a_{i,j}),
$$
\n(2.5)

is based on a family of probability density functions  $(p.d.f.)$   $f(y;a)$  with a nonnegative y and a non-negative scalar parameter a. For the Poisson and exponential p.d.f., the model in (2.5) simplifies, leading to a generalized linear model (GLIM).

Static models with network matrices can be also applied to data on a set of individuals  $i = 1, \ldots, L$  where at time t the number  $y_{i,j,t}$  of messages sent by individual i to individual j are observed. Alternatively, the i.i.d. observations  $y_t$  can be matrices containing the investments of bank  $i$  in industrial sector  $j$  at time  $t$ , gravity matrices summarizing international trades of countries [Chen et al. (2021), Section 6], or matrices of numbers of stocks of firms head-quartered in city  $j, j = 1, \ldots, m$  held by mutual fund manager  $i, i = 1, \ldots, n$  at time t.

#### 2.2.2 Panel of individual qualitative histories.

Let us consider a panel of L individuals,  $l = 1, \ldots, L$ . Each individual is characterized by a qualitative state  $i = 1, \ldots, n$ , that can be observed at any time t. The qualitative individual histories can be quantified and represented by  $n$ -dimensional vectors  $(Y_{l,t}, t = 1, \ldots, T)$ ,  $l = 1, \ldots, L$ , where  $Y_{l,t}$  has  $\{0, 1\}$  entries that sum up to 1.

The dynamic model can be defined by assuming that :

i) the individual histories are independent from one another; ii) each individual history corresponds to a Markov chain with transition matrix A.

The above assumptions imply that the population of interest is homogeneous. Under these assumptions, the individual histories can be aggregated without a loss of information and replaced by the counts of individuals in each state :

$$
Y_t = \sum_{l=1}^{L} Y_{l,t}.
$$
\n(2.6)

Then, the sequence of multivariate counts  $Y_t$  is also a Markov process with the conditional p.d.f. obtained from the convoluate of different multinomial distributions :

$$
l(y_t|y_{t-1};A) = \Pi_{i=1}^n l_i(y_{i,t}|y_{i,t-1};a_i),
$$
\n(2.7)

where  $l_i$  denotes the p.d.f. of the multinomial distribution  $M(y_{i,t-1}; a_i)$  and  $a_i$  denotes the  $i^{th}$ -row of the transition matrix A.

In this example, individuals are replaced by homogeneous segments that diminishes the dimensionality of the network. For example, the corporates can be replaced by industrial sectors, households by age categories, or animals by species [Donnet and Robin (2021) and the references therein. Then, the diagonal elements of matrix  $A$ depict the interactions within a segment and the off-diagonal elements of A represent the interactions between the segments. In general, the diagonal elements  $a_{ii}$  are not equal to zero.

#### 2.2.3 Dynamic model for a non-negative random vector

Let us consider a parametric model for a non-negative random vector  $l(y; \theta)$ , where y and the parameter vector  $\theta$  have the same dimension n, and both vectors y and  $\theta$ are non-negative. Then, a dynamic model for  $(y_t)$  can be defined as :

$$
l(y_t|y_{t-1};A) = l(y_t; Ay_{t-1}),
$$
\n(2.8)

where the contagion matrix A is non-negative and of dimension  $(n, n)$ . This model can be easily extended to include an intercept  $\mu$ , i.e.  $l(y_t|y_{t-1};A) = l(y_t, Ay_{t-1} + \mu)$ .

i) When this parametric model represents the dynamics of  $n$  independent Poisson variables, we obtain a multivariate Autoregressive Conditional Poisson (ACP) model with :

$$
l(y_t|y_{t-1};A) = \Pi_{i=1}^n \left[ \frac{1}{y_{it!}} \exp(-a_i y_{t-1})(a_i y_{t-1})^{y_{it}} \right]
$$
  

$$
= \Pi_{i=1}^n \left( \frac{1}{y_{it!}} (a_i y_{t-1})^{y_{it}} \right) \exp(-e' A y_{t-1}), \qquad (2.9)
$$

where  $a_i$  is the  $i^{th}$  row of matrix A and e is a vector with unitary elements [Cameron and Trivedi (2008), Section 7.5, eq. 7.42, Fokianos (2022)] .

This specification differs from an exponential specification:  $y_{it}|y_{t-1} \sim \mathcal{P}[\exp(a_i y_{t-1})],$ considered in Chen et al. (2021), Example 3, which does not require matrix A to be non-negative. However, while the dynamic of model (2.9), possibly with an intercept, is compatible with the stationarity of process  $(y_t)$ , the above approach of Chen et al. (2021) leads to explosive trajectories, due to the exponential transform of a.

ii) When the parametric model represents the dynamics of  $n$  independent exponential variables, we get a multivariate exponential autoregressive model:

$$
l(y_t|y_{t-1};A) = \Pi_{i=1}^n [(a_iy_{t-1}) \exp(-a_iy_{t-1}y_{it})].
$$
\n(2.10)

This dynamic model can be used to study joint evolutions of the gross domestic product in a set of  $L$  countries. Then, matrix  $A$  is unobserved and needs be estimated under some mild constraints to provide a proxy of the international trading network. It can also be used for joint analysis of observed (implied or realized) daily volatilities, with the individual index  $i$  assigned either to the stocks, or hours of a trading day, in the spirit of multiplicative volatility models [Engle and Rangel (2008), Hafner and Linton (2010)].

**Remark 2**: The dynamic model (2.8) can be extended to  $l(y_t|y_{t-1};A) = l(y_t;Az_{t-1}),$ where  $z_{t-1}$  is a nonnegative vector function of  $y_{t-1}$ . Such a transformation appears in structural models used in epidemiology and other applications such as the analysis of cyberattacks [Fahrenwaldt et al. (2018), Lu et al. (2023)], adoption of a new technologies [Brock and Durlauf (2010), Blume et al. (2011)] and the Susceptible-Infected-Recovery (SIR) model with multiple transmissions, where the components of y are the counts of infected individuals [resp. of the new adoptions of the product] in different segments of the populations. In the SIR model,  $z$  is a quadratic function of y [see e.g. Gourieroux and Lu  $(2023)$ ].

#### 2.2.4 Contagion of defaults

The structural models of corporate defaults represent defaults arising when the liability of a firm falls below its asset value. When only defaults are observed, the individual assets and liabilities become latent non-negative variables which can be represented by a dynamic network model (2.8):

$$
\tilde{l}(z_t|z_{t-1};A) = \tilde{l}(z_t;Az_{t-1}),
$$
\n(2.11)

where  $z_t$  has dimension 2L and L is the number of firms. The elements of  $z_t$  represent the assets and liabilities of firms. For example,  $z_{1lt}$ ,  $z_{2lt}$  are the characteristics of firm l at date t. The observed variables are defaults, represented by a vector  $y_t$  of dimension L and defined by :

$$
y_{lt} = 1
$$
, if  $z_{1lt} > z_{2lt} = 0$ , otherwise. (2.12)

Equations  $(2.11)$  and  $(2.12)$  can be interpreted as a nonlinear state space model. From the state equation (2.11) and measurement equation (2.12), it follows that the transition density of observed variables is  $l(y_t|y_{t-1};A)$ , involving the complete histories of defaults of all firms.

All dynamic models given above have nonlinear dynamics due to the non-negativity constraints on the values of observed variables. They are also nonlinear with respect to parameter A. These nonlinear features distinguish the class of models considered in this paper from the major part of literature on networks in econometrics. "Identification and measurement of network phenomena has drawn attention in fields as diverse as macroeconomics, finance not discussed in the review" [De Paula (2017)]. Other fields with nonlinear dynamic models are also left out of the recent survey on partial identification by Kline and Tamer (2023). The nonlinear dynamic models considered in our paper are suitable for applications to those other fields, including ecological economy, monetary economy, cyber risk and epidemiology, in addition to finance.

The identification in this extended class of models eliminates the use of powers of matrix A. Indeed, matrices  $A, A^2, A^3$  no longer play a role, as the conditional distribution of  $y_t$  given  $y_{t-h}$  does not necessarily depend on parameter A through  $A^h$ only [see de Paula (2017), p. 272 for a discussion]. In fact, under the Markov of order 1 assumption, all information on  $A$  is contained in the first lag of variables in the likelihood function. Moreover, the response of a nonlinear dynamic system to a shock at date  $t$  depends on the current environment of  $y_t$ . The size of a shock effect depends on the environment, so that even a small shock can have a large impact. This has important consequences on how the small-valued elements  $a_{ij}$  of matrix A should be treated. More precisely, even if  $a_{ij}$  is small, it should not be set equal to zero arbitrarily, by applying an automatic LASSO penalty. Especially, when the dynamic system is close to a tipping point, a small connection can trigger a significant change in the system, as evidenced by the histories of corporate defaults, chains of business failures, or inter-bank liquidity shortages.

## 2.3 Latent Heterogeneity and Ranking

Taking into account the non-negativity condition can be beneficial for the interpretation and visualization of networks. Let us now introduce a normalized NMF, which is suitable for structural interpretations.

#### 2.3.1 Alternative parametrization

The NMF (2.4) can be normalized and written alternatively as :

$$
A = a \sum_{k=1}^{K} \pi_k \beta_k^* \gamma_k^{*'},
$$
\n(2.13)

where  $a = e'Ae = \sum_i \sum_j a_{i,j} > 0, \pi_k \geq 0, k = 1, \ldots, K$ , with  $\sum$ K  $k=1$  $\pi_k = 1, \beta_k^* \geq 0, \gamma_k^* \geq$  $0, k = 1, ..., K$ , with :  $\beta_k^{*'}$  $k^{*'}e = \gamma_k^{*'}$  $k^{\ast}e = 1, k = 1, \ldots, K.$ 

In this decomposition  $\pi = (\pi_1, \ldots, \pi_K)'$ ,  $\beta_k^*, \gamma_k^*, k = 1, \ldots, K$  can be interpreted as discrete probability distributions. More precisely, the normalised matrix  $A/a$  can be interpreted as a joint probability distribution, and its decomposition  $\sum$ K  $k=1$  $\pi_k \beta_k^* \gamma_k^{*'}$  $\kappa$ <sup>\*'</sup> as a mixture of independent joint distributions. In this respect, our analysis is related to the literature on partial identification of finite mixtures [see Hall and Zhou (2003), Kasahara and Shimotsu (2009), Henry et al. (2014), and Section 5. When  $K = 1$ , the NMF becomes :  $A = a\beta_1^* \gamma_1^{*'}$  $\frac{\ast'}{1}$ .

The representation (2.13) solves some identification issues due to the identification of factorial directions  $\beta_k, \gamma_k$  up to positive scalars. It also allows us to bound the set of NMF's written under this form.

#### 2.3.2 Rankings

To motivate the alternative parametrization (2.13), let us consider a dynamic model of count variables in epidemiology, such as the dynamic contagion model (2.8), where the components of  $y_t$  are counts of infected individuals in  $L$  homogenous segments of the population. In addition, let us assume that <sup>[4](#page-0-0)</sup>:

$$
E_{t-1}y_t = Ay_{t-1}.\tag{2.14}
$$

If  $K = 1$ , we get :  $E_{t-1}y_t = a\beta_1^*\gamma_1^{*'}$  $j_1^{*'}y_{t-1}$ , or equivalently :

$$
E_{t-1}y_{i,t} = \sum_{j=1}^{L} a_{ij}y_{j,t-1} = a\beta_{1i}^{*} \sum_{j=1}^{L} \gamma_{1j}^{*}y_{j,t-1}.
$$
 (2.15)

The contagion parameters  $a_{ij}$  can be decomposed as:  $a_{ij} = a\beta_{1i}^*\gamma_{1j}^*$ , where a is a global contagion effect,  $\beta_{1i}^*$  an index of vulnerability of segment i to the infection and  $\gamma_{1j}^*$  a measure of viral load of segment j<sup>[5](#page-0-0)</sup>. Therefore the segments  $i = 1, \ldots, L$  can be ranked with respect to their vulnerability  $\beta_{1i}^*$  and their infectiosity  $\gamma_{1i}^*$ .

When K is larger than or equal to 2, a latent heterogeneity of segments arises with heterogeneity distribution  $\pi$ . Then, the segments can be ranked with respect to their vulnerability ratings  $\beta_{ki}^*, k = 1, ..., K$ , and infectiosity ratings  $\gamma_{k,i}^*, k = 1, ..., K$ . The potential interpretations of parameters  $\beta^*, \gamma^*, \pi$  depend on the application of interest. In the analysis of internet diffusion of messages (see Section 2.2.1), parameters

<sup>&</sup>lt;sup>4</sup>This constraint on the conditional mean implies  $y_t = Ay_{t-1} + u_t$ , where  $u_t$  is a martingale difference sequence with  $E_{t-1}(u_t) = 0$ . It does not imply a linear dynamic model :  $y_t = Ay_{t-1} + u_t$ , where  $u_t$  is a strong (i.i.d.) white noise. Hence, this linear model is not suitable for the analysis of nonlinear shock effects, i.e. nonlinear impulse response functions.

<sup>&</sup>lt;sup>5</sup>The decomposition of  $a_{ij}$  is multiplicative, while an additive decomposition:  $a_{ij} = \tilde{\alpha} + \tilde{\beta}_i + \tilde{\gamma}_j$ , is used in the panel literature with two-ways fixed effects [see e.g. the running examples in Fernandez-Val and Weidner (2016) for probit and Poisson models]. In our framework it is not possible to transform the multiplicative form into an additive one by taking the logarithms of  $a_{ij}, \beta_{1i}^*, \gamma_{1j}^*$ , as there can exist individuals with zero value of  $\beta_1^*$ , such as the vaccinated or naturally immunized individuals. From the identification perspective, these zero values can be informative and should not be disregarded. Moreover, as noted in Chen et al. (2021), such interactive effects can capture network features, such as homophily and clustering.

 $\beta^*$  (resp.  $\gamma^*$ ) can be used for ranking the receivers and senders, or followers and influencers, in trade networks for ranking the importers and exporters, in citation networks for ranking the citees and citors, in cyber risk networks for ranking of firms by vulnerability and hackers by their impact.

## 3 Identification of the Nonnegative Factorization

In general, the true NMF is not point-identified. This section discusses the identification issues and derives the identified set of NMFs for a given non-negative rank.

## 3.1 The general framework

The parametric model depends on the non-negative matrix factorization :

$$
A = BC' = (\beta_1, \dots, \beta_K)(\gamma_1, \dots, \gamma_K)' = \sum_{k=1}^K \beta_k \gamma'_k,
$$

where K denotes the non-negative rank and  $\beta_k, \gamma'_k s$  are the factorial directions. In practice, the structural parameters are  $K, \beta_k, \gamma_k, k = 1, ..., K$ ,. There exists a large body of literature on the lack of identification of these parameters for a given matrix A. It is easy to see that the factorization is not unique because the same matrix A is obtained from a permutation of index  $k$  and rescaling by a positive scalar, i.e. by replacing  $\beta_k, \gamma_k$  by  $\sigma_k \beta_k$ , and  $\gamma_k / \sigma_k$ , respectively, for a positive scalar  $\sigma_k$ . This identification issue is easily solved by a normalization. The more complicated identification issues arise when  $K \geq 2$  and are discussed below under the following assumption :

**Assumption A.3**: The non-negative rank of matrix A is equal to the rank of A.

This assumption is not very stringent, even though some examples of nonnegative matrices with  $Rk_{+}(A) > Rk(A)$  have been given in the literature. It is useful to describe and parametrize the set of admissible NMF's when  $Rk_{+}(A) = Rk(A)$ . Assumption A.3 can be written under the following equivalent forms :

Assumption A.3 is satisfied  $\Leftrightarrow$  the vectors  $\beta_1, \ldots, \beta_K$  are linearly independent and the vectors  $\gamma_1, \ldots, \gamma_K$  are linearly independent  $\Leftrightarrow B'B$  and  $C'C$  are invertible<sup>[6](#page-0-0)</sup>.

<sup>6</sup>The NMF representation is different from the Singular Value Decomposition (SVD) of matrix

Assumption A.3 implies that the matrices  $\beta_k \gamma_l', k, l = 1, ..., K$  are linearly independent as shown in Lemma 1 below :

**Lemma 1**: Under Assumption A.3,  $B\Delta C' = 0 \Rightarrow \Delta = 0$ .

**Proof :** Indeed  $B\Delta C' = 0$  implies  $B'B\Delta C'C = 0$ , and then  $\Delta = 0$ , since  $B'B$ and  $C'C$  are invertible. QED

Under assumptions A.1-A.3, the non-negative rank and the rank of  $A_0$  are known. Let us consider another factorization of matrix A without taking into account the non-negativity conditions on  $\beta_k, \gamma'_k$  s. Since the range of A (resp. A') is the space spanned by  $\beta_1, \ldots, \beta_K$  (resp. by  $\gamma_1, \ldots, \gamma_K$ ), an alternative factorization is :  $A =$ B G H'C' =  $\tilde{B}\tilde{C}'$ , where  $\tilde{B} = B G, \tilde{C} = C H$  and G, H are invertible matrices <sup>[7](#page-0-0)</sup>. Moreover, we have : B G H'  $C' = B C'$ , and by Lemma 1 we deduce that  $H' = G^{-1}$ . Therefore we get :  $\tilde{B} = BG, \tilde{C} = C(G')^{-1}.$ 

In addition, because of the definition of factors up to the permutation and (signed) scale effects, we can choose  $G$  of the type :

$$
G = Q \text{ diag } \sigma,\tag{3.1}
$$

where  $\sigma = (\sigma_k)$ ,  $\sigma_k > 0$ , and Q a  $(K, K)$  invertible matrix with diagonal elements equal to 1. Then,  $(G')^{-1} = (Q')^{-1}$  diag  $(1/\sigma)$ . By taking into account the nonnegativity conditions, we get a constructive characterization of the identified set :

**Proposition 1 :** For a specific factorization of matrix  $A_0$  :  $A_0 = B_0 C'_0, B_0 \ge$  $0, C_0 \geq 0$  (referred to as the origin, henceforth) of the identified set), all observationally equivalent non-negative factorizations are such that :

$$
\tilde{B}=B_0Q \text{ diag }\sigma, \ \tilde{C}=C_0(Q')^{-1} \text{ diag }(1/\sigma), \ \sigma_k>0, \forall k=1,\ldots,K,
$$

where the matrix  $Q$  is invertible, with unitary diagonal elements, and such that :

$$
B_0 Q \geq 0, C_0 (Q')^{-1} \geq 0.
$$

A. In SVD the identification issue is solved by introducing the orthonormality restriction  $B'B =$  $C'C = Id.$  Orthogonality is not possible in our framework since  $\beta'_k \beta_l$  is always non-negative. The orthogonality condition would imply  $\beta_{ik}\beta_{il} = 0$ ,  $\forall i$ , and contradict assumption A.1 [see the discussion in Online Appendix 4].

<sup>&</sup>lt;sup>7</sup>The columns of  $\tilde{B}$  and the columns of  $\tilde{B}$  are two bases of the range of  $AA'$  equal to the range of A. Therefore they are in a one-to-one relationship represented by an invertible matrix G.

The non-negative factorization is said to be essentially unique, or simply unique [see Laurberg et al. (2008)], if  $Q = Id$  is the only solution to the set of inequalities given in Proposition 1.

Proposition 1 shows that the identified set is fully parametrized once an origin is chosen. The parametrization in  $Q$  is not linear. Then, in general, the identified set is neither convex, nor even star-convex.

### 3.2 Special Cases

#### **3.2.1** Case  $K = Rk_{+}(A) = 1$

When  $K = 1, A = \beta_1 \gamma'_1, Q = (q_{11}) = (1)$ , the non-negative factorization is essentially unique, i.e. the NMF is (essentially) point identified. Moreover, we note that :

$$
A\beta_1 = \beta_1(\gamma'_1\beta_1), A'\gamma_1 = \gamma_1(\beta'_1\gamma_1).
$$

It follows that  $\beta_1$  (resp.  $\gamma_1$ ) is an eigenvector of A (resp. A') associated with the eigenvalue  $\gamma'_1 \beta_1 = \beta'_1 \gamma_1$ , which is strictly positive. By the Perron-Froebenius Theorem [see Meyer  $(2000)$ ] the non-negative matrix A (resp. A') has a unique eigenspace of dimension 1 generated by a non-negative eigenvector, here  $\beta_1$  (resp.  $\gamma_1$ )<sup>[8](#page-0-0)</sup>.

### 3.2.2 Case  $K = Rk_{+}(A) = 2$

For  $K = 2$ , we get  $Q =$  $\begin{pmatrix} 1 & q_{12} \end{pmatrix}$  $q_{21}$  1  $\bigg),$  and  $(Q')^{-1} = \frac{1}{1 - q_{12}q_{21}} \begin{pmatrix} 1 & -q_{21} \\ -q_{12} & 1 \end{pmatrix}$  $-q_{12}$  1  $\setminus$ . Without the subscript 0 for the "specific factorization of the true matrix", the inequality restrictions in Proposition 1 become :

$$
\begin{cases} \n\beta_{1,j} + q_{21}\beta_{2,j} \ge 0, \ q_{12}\beta_{1,j} + \beta_{2,j} \ge 0, \ j = 1, \dots, n, \\
\frac{1}{1 - q_{12}q_{21}}(\gamma_{1,j} - q_{12}\gamma_{2,j}) \ge 0, \ \frac{1}{1 - q_{12}q_{12}}(-q_{21}\gamma_{1,j} + \gamma_{2,j}) \ge 0, j = 1, \dots, m.\n\end{cases}
$$

It is easy to check that these inequalities imply  $1 - q_{12}q_{21} > 0$ . The latter inequality is satisfied, in particular, when we consider local identification in a neighbourhood of

<sup>&</sup>lt;sup>8</sup>It is easy to check that  $\beta_1$  (resp.  $\gamma_1$ ) is also an eigenvector of AA' (resp. A'A). Therefore they are both elements of a singular value decomposition (SVD) (see Remark 3 in Section 4.1). Such a model is commonly applied in the network literature with an adjacency matrix, whose leading left and right eigenvectors of the SVD are used to define the "so-called" hub and authority centralities, respectively [see Cai et al. (2021) and the references therein].

 $B = (\beta_1, \beta_2), C = (\gamma_1, \gamma_2)$ , i.e. in a neighbourhood of  $q_{12} = q_{21} = 0$ . Then, the system of inequalities is equivalent to :

$$
\begin{cases} \n\beta_{1,j} + q_{21}\beta_{2,j} \ge 0, \forall j, \text{ with } \beta_{2,j} > 0, \ -q_{21}\gamma_{1,j} + \gamma_{2,j} \ge 0, \forall j, \text{ with } \gamma_{1,j} > 0, \\
q_{12}\beta_{1,j} + \beta_{2,j} \ge 0, \ \forall j, \text{ with } \beta_{1,j} > 0, \ \gamma_{1,j} - q_{12}\gamma_{2,j} \ge 0, \forall j, \text{ with } \gamma_{2,j} > 0,\n\end{cases}
$$

or

$$
\begin{cases} q_{12} \ge \sup_{j:\beta_{1,j}>0} (-\beta_{2,j}/\beta_{1,j}), q_{12} \le \inf_{j:\gamma_{2,j}>0} (\gamma_{1,j}/\gamma_{2,j}), \\ q_{21} \ge \sup_{j:\beta_{2,j}>0} (-\beta_{1,j}/\beta_{2,j}), q_{21} \le \inf_{j:\gamma_{1,j}>0} (\gamma_{2,j}/\gamma_{1,j}). \end{cases}
$$

From Proposition 1, we deduce the following result :

**Proposition 2 :** For  $K = Rk_{+}(A) = 2$ , the admissible matrices  $Q =$  $\begin{pmatrix} 1 & q_{12} \end{pmatrix}$  $q_{21}$  1  $\setminus$ are such that :

$$
- \inf_{j:\beta_{2,j}>0} (\beta_{1,j}/\beta_{2,j}) \le q_{21} \le \inf_{j:\gamma_{1,j}>0} (\gamma_{2,j}/\gamma_{1,j}),
$$
  

$$
- \inf_{j:\beta_{1,j}>0} (\beta_{2,j}/\beta_{1,j}) \le q_{12} \le \inf_{j:\gamma_{2,j}>0} (\gamma_{1,j}/\gamma_{2,j}).
$$

Therefore, among the  $2(n + m)$  inequality restrictions in Proposition 1, only four are active and the remaining ones are redundant. We deduce the necessary and sufficient exclusion conditions for essential uniqueness [Brie (2015)].

**Corollary 1 :** Under Assumption A.3 and for  $K = Rk_{+}(A) = 2$ , the nonnegative factorization is essentially unique if and only if there exists at least one index  $j_1$  such that  $\beta_{1,j_1} = 0, \beta_{2,j_1} > 0$ , one index  $j_2$  such that  $\beta_{1,j_2} > 0, \beta_{2,j_2} = 0$ , one index  $j_3$  such that  $\gamma_{1,j_3} = 0, \gamma_{2,j_3} > 0$  and one index  $j_4$  such that  $\gamma_{1,j_4} > 0, \gamma_{2,j_4} = 0$ .

Without the "exclusion" restrictions of Corollary 1 and under Assumption A.2 iii), there exists a multiplicity of admissible non-negative factorizations that are easily deduced from one of them. This is the case we are interested in.

Let us now discuss the degree of under-identification. Since the identification issue of factorial directions up to multiplicative non-negative scalars is solved in representation (2.13), we can focus our attention on the identification of  $\pi_k, \beta_k^*, \gamma_k^*, k = 1, 2$ . For K=2, the decomposition (2.13) is:  $A = a[\tilde{\pi}_1 \tilde{\beta}_1^* \tilde{\gamma}_1^{*'} + \tilde{\pi}_2 \tilde{\beta}_2^* \tilde{\gamma}_2^{*'}]$  $\binom{*'}{2}$ .

It is easy to derive the closed-form parametrization of the identified set for the given origin  $(\beta_1^*, \beta_2^*), (\gamma_1^*, \gamma_2^*)$  (see online Appendix 1):

$$
\tilde{\beta}_1^* = p_1 \beta_1^* + (1 - p_1) \beta_2^*, \text{ with } p_1 = \beta_1' e / (\beta_1' e + q_{21} \beta_2' e),
$$
  
\n
$$
\tilde{\beta}_2^* = p_2 \beta_1^* + (1 - p_2) \beta_2^*, \text{ with } p_2 = q_{12} \beta_1' e / (q_{12} \beta_1' e + \beta_2' e),
$$
  
\n
$$
\tilde{\gamma}_1^* = p_3 \gamma_1^* + (1 - p_3) \gamma_2^*, \text{ with } p_3 = \gamma_1' e / (\gamma_1' e - q_{12} \gamma_2' e),
$$
  
\n
$$
\tilde{\gamma}_2^* = p_4 \gamma_1^* + (1 - p_4) \gamma_2^*, \text{ with } p_4 = -q_{21} \gamma_1' e / (q_{21} \gamma_1' e - \gamma_2' e),
$$
  
\n
$$
\tilde{\pi}_1 / \tilde{\pi}_2 = (\beta_1' e + q_{21} \beta_2' e) (\gamma_1' e - q_{12} \gamma_2' e) / (q_{12} \beta_1' e + \beta_2' e) (-q_{21} \gamma_1' e + \gamma_2' e).
$$

**Corollary 2**: For  $Rk(A) = Rk_{+}(A) = 2$ , the components  $\pi_k, \beta_k^*, \gamma_k^*, k = 1, 2$  are not identifiable, with a degree of under-identification equal to 2.

The set of admissible decompositions (2.13) is described in closed form by means of additional parameters  $q_{12}, q_{21}$ . The conditions on  $q_{12}, q_{21}$  derived in Proposition 2 remain valid when  $\beta$ ,  $\gamma$  are replaced by  $\beta^*, \gamma^*$ .

#### 3.2.3 General Case

For  $K = 2$ , the identified set is described by two parameters  $q_{12}, q_{21}$  satisfying  $2(m+n)$ inequality restrictions. These restrictions are linear and Proposition 2 shows that only 4 of them are active. In the general case, the identified set is parametrized by  $K(K-1)$  parameters, which are the off-diagonal elements of matrix Q. Therefore, the identified set is a manifold of fixed dimension. These parameters satisfy  $K(n+m)$ inequality restrictions by Proposition 1. The degree of under-identification increases rather quickly with the non-negative rank and the subset of active restrictions cannot be derived analytically. To determine these restrictions and provide a simplified definition of the identified set, numerical algorithms are needed, such as the Active Set Sequential Quadratic Programming (SQP) algorithm, some of them being available from Artelys Knitro [see Gill et al. (2002), Liu (2005)] at www.artelys.com.

A similar problem arises in sharp set identification for discrete choice models and treatment effects analysis. The main difference is that those inequality restrictions

are linear, e.g. as for  $K = 2$  in Section 3.2.2<sup>[9](#page-0-0)</sup>, while in our framework they are nonlinear when  $K \geq 3$ . Indeed, for  $K = 3$ , we get :  $Q =$  $\sqrt{ }$  $\mathcal{L}$ 1 q<sup>12</sup> q<sup>13</sup>  $q_{21}$  1  $q_{23}$  $q_{31}$   $q_{32}$  1  $\setminus$  and, up to the determinant, matrix  $(Q')^{-1}$  has cofactor elements, such as  $1 - q_{23}q_{32}$  for instance, that are quadratic in Q (more generally, they are polynomials of degree less or equal to  $K-1$ ). Moreover, it is easy to see that the identified set is neither convex, nor star-convex, in general.

## 4 Statistical Inference

A major challenge for statistical inference is the lack of identification of the true NMF. This identification issue can be addressed either by introducing identification restrictions and applying the standard maximum likelihood approach, or by estimating the set of all identifiable NMF directly, either pointwise, or from a confidence set<sup>[10](#page-0-0)</sup>. These two approaches are linked. For example, the set of all identifiable NMF's can be deduced from one of them, as shown in Section 3.3.2 for  $Rk_{+}(A) = 2$ . Therefore, we can first identify one NMF (i.e. an origin) and next deduce all the remaining ones from the origin. This will provide an estimator of the identified set. Alternatively, we can derive directly an asymptotic confidence set of the identified set at a given confidence level by inverting appropriately a test procedure.

In this respect, the identified set described in Proposition 1 has the general form<sup>[11](#page-0-0)</sup> considered in Shi and Shum (2015), eq. (1.1)-(1.2). However, their results can only be used if the origin is well-defined and has a consistent and asymptotically normally distributed estimator. These are the main points considered in this section.

The proposed method proceeds as follows: 1) We show the convergence of the set of maximum likelihood estimators of  $B, C$  to the identified set; 2) For a fixed number of observations, a multiplicity of ML estimators is obtained. Therefore, we introduce an alternating ML algorithm to fix the selected ML estimator for any T. This sequence of alternating ML estimators is well-defined and converges to the identified set, but

<sup>&</sup>lt;sup>9</sup>See the linear moment functions defining the inequalities in Chernozhukov et al. (2007), Section 2.1.

<sup>&</sup>lt;sup>10</sup>We distinguish between the pointwise estimated identified set and the confidence set of the identified set at a given confidence level.

<sup>11</sup>up to the introduction of nuisance slackness parameters [Shi and Shum (2015), Remark p. 497].

not necessarily pointwise to a given element of this set; 3) Next, we fix an origin in the interior of the identified set, obtained as the solution of an auxiliary optimization. The objective function of that auxiliary optimization is optimized with respect to an alternative parametrization in  $(a, \pi, \beta^*, \gamma^*)$ . This step provides a consistent estimator of an origin in the identified set.

The new Identifying Maximum Likelihood (IML) approach is used to estimate the identified set, analyze the key properties of this set estimator and derive its asymptotic distributional properties. We study analytically the properties of the distribution of the identified set estimator by using the random set theory in [Molchanov, Molinari  $(2018)$ . In addition, we derive the ML estimator of matrix A under a given nonnegative rank constraint and the efficiency bound of this identifiable matrix-parameter estimator. The approach is feasible because of the properties of the maximum likelihood estimator reviewed below.

## 4.1 The ML approach

In our framework of nonlinear models with partial identification, the ML estimator(s) does not have a closed-form expression. Hence, we have to distinguish the properties of the statistical convergence of ML estimators to the true identified set when the number T of observations tends to infinity from the numerical convergence of the algorithm used to approximate the ML estimator when the number of iterations tends to infinity<sup>[12](#page-0-0)</sup>.

#### 4.1.1 Consistency of ML Estimators

Let us assume that the network model is well-specified and the true transition is :

$$
l(y_t|y_{t-1};A_0) = l(y_t|y_{t-1};B_0C'_0),
$$
\n(4.1)

with a non-negative rank  $K_0$  of matrix  $A_0$  assumed to be known. We consider a

<sup>&</sup>lt;sup>12</sup>In this respect, our analysis differs from the literature on exact NMF that implicitly assumes that  $A_0$  is observed and examines the computational complexity, i.e. the possibility of finding a decomposition of  $A_0$  in polynomial time, called the NP-hardness [see Gillis (2020) for a survey]. This literature disregards the uncertainty of observations.

constrained ML maximization, providing:

$$
(\hat{B}_T, \hat{C}_T) = \arg\max_{B \ge 0, C \ge 0} \sum_{t=1}^T \log l(y_t | y_{t-1}; BC').
$$
 (4.2)

Due to the identification issue, there exists a large multiplicity of solutions to the finite sample optimization (4.2). Under additional regularity conditions given in online Appendix 2, the above set of ML estimators is consistent.

Proposition 3 : Under Assumptions A.1-A.3 and a.1 in Appendix 2, the set of ML estimators  $(\hat{B}_T, \hat{C}_T)$  converges to the set  $\mathcal{A}_0$  of NMF associated with  $A_0 = B_0 C'_0$ , when  $T$  tends to infinity.

More precisely, let  $d(.,.)$  denote the Euclidean distance on  $\mathbb{R}^{K(n+m)}$ ,  $NMF_0 =$  $\{(B, C), B \ge 0, C \ge 0, \text{ with } BC' = B_0 C_0'\}, \text{ and } \mathcal{D}[(\hat{B}_T, \hat{C}_T), NMF_0] = \min_{(B, C) \in NMF_0}$  $d[(\hat{B}_T, \hat{C}_T), (B, C)]$ , then  $\mathcal{D}[(\hat{B}_T, \hat{C}_T), NMF_0]$  tends to zero, when T tends to infinity. Under the regularity conditions, this convergence is uniform in  $(\hat{B}_T, \hat{C}_T)$ . It will also imply the convergence to this set of the well-defined alternating ML estimator introduced in the next section. Thus,  $(\hat{B}_T, \hat{C}_T)$  does not necessarily converge to the true factorization  $(B_0, C_0)$  due to the identification issue, but for a large T,  $(\hat{B}_T, \hat{C}_T)$  is close to another admissible NMF that can depend on T.

#### 4.1.2 Alternating ML (AML) Algorithm

A ML estimator of matrix  $A = BC'$  does not always have a closed-form. In practice, it is computed numerically from an algorithm, such as a Newton-Raphson algorithm. In our framework of partial identification, a Newton-Raphson type of algorithm cannot be used jointly for  $B$  and  $C$ . The reason is that each iteration requires an inversion of the Hessian matrix, which is not invertible due to the identification issue. The AML algorithm, also called the block mirror descent (BMD) algorithm [Hien and Gillis (2021), or zig-zag algorithm see Hautsch at al. (2023) solves this issue. In the presence of a multiplicity of NMFs, we apply an AML algorithm [see, e.g. Gourieroux, Monfort and Renault (1990), Kim and Park (2008), Hastie et al. (2015), Gillis (2020), Chapter 8 for alternating least squares, Hien and Gillis (2021) for NMF]. We observe that, even if the factorization  $B, C$  is not identifiable,  $B$  (resp. C) is identi-

fiable when  $C$  is known (resp.  $B$  is known). This leads to the following alternating ML algorithm, where at step p,  $\hat{B}_{p,T}, \hat{C}_{p,T}$  is computed and then  $\hat{B}_{p+1,T}, \hat{C}_{p+1,T}$  are recursively obtained as follows:

$$
\hat{B}_{p+1,T} = \arg \max_{B \ge 0} \sum_{t=1}^{T} \log l(y_t | y_{t-1}; B \hat{C}'_{p,T}), \tag{4.3}
$$

$$
\hat{C}_{p+1,T} = \arg \max_{C \ge 0} \sum_{t=1}^{T} \log l(y_t | y_{t-1}; \hat{B}_{p+1,T} C'). \tag{4.4}
$$

By construction, this algorithm produces at each iteration  $p$  a higher value of the log-likelihood function than at iteration  $p-1$ .

We have to distinguish the ML estimator from the alternating ML estimator obtained from the algorithm  $(4.3)-(4.4)$ . As mentioned earlier, when some parameters are not identifiable there is a multiplicity of ML estimators of  $B$  and  $C$ . However, there is a unique sequence of alternating ML estimators for the given starting values, even though the AML algorithm does not necessarily numerically converge pointwise, due to the identification issue. Moreover, even if it converges numerically to a global  $\text{maximum}^{13}$  $\text{maximum}^{13}$  $\text{maximum}^{13}$ , the limit can depend on the starting value and does not necessarily correspond to a point in the interior of the identified set. The dependence on the initialization has been observed in practice and led to a literature on the optimal choice of the starting value [see the discussion in Gillis (2020), Chapter 4 and Esposito (2021)]. This shows a confusion between the effect of non-identification and the fact that the algorithm can stop at a local, instead of global maximum in some cases.

The discussion of starting values is different when the focus is on the estimation of matrix A under the NMF restriction. Indeed, the matrix A is identifiable even if the pair of matrices  $B$  and  $C$  is not. For this identifiable parameter matrix  $A$  it is possible to prove that the AML estimator and the infeasible ML estimator have identical asymptotic distributions, if the starting value of the AML algorithm is a consistent estimator of A [see Vrahalis et al. (2003), Hautsch et al. (2023)]. Such starting values could be derived from an unconstrained estimator of A if the number  $n^2$  of elements in A is not too large, compared to the number of observations  $nT$ .

<sup>13</sup>See, Grippo and Sciandrosse (2000) for the numerical convergence conditions.

# 4.2 Nonnegative rank  $K_0 = Rk_+(A_0) = 1$

As pointed out in Section 3.2.1, the NMF is essentially unique for  $K_0 = 1$ . The assumption  $K_0 = 1$  greatly simplifies the estimation and explains its frequent use in applied econometrics [see e.g. Cai et al. (2021)]. Because the ML estimator is unique in this case, it can be computed from a standard Newton-Raphson algorithm.

We introduce the identification restrictions:  $A = a\beta\gamma'$ , with  $a > 0, \beta \geq 0, \gamma \geq 0$ , and  $\beta' e = \gamma' e = 1$ . Let us now consider a dynamic model described in Section 2.2.3, which is the latent parametric model  $l(y; \theta)$ , where  $\theta$  is replaced by  $Ay_{t-1}$  (see, Section 2.2.3). We get :

$$
l(y_t|y_{t-1};A) = l(y_t; a\beta\gamma'y_{t-1}).
$$
\n(4.5)

The partial derivatives of the log-likelihood with respect to  $a, \beta, \gamma$  are easily derived from the partial derivatives of the latent log-likelihood with respect to  $\theta$ . We have :

$$
\begin{cases}\n\frac{\partial \log l}{\partial \beta}(y_t|y_{t-1};A) = a\gamma' y_{t-1} \frac{\partial \log l}{\partial \theta}(y_t; a\beta \gamma' y_{t-1}), \\
\frac{\partial \log l}{\partial \gamma}(y_t|y_{t-1};A) = ay_{t-1} \beta' \frac{\partial \log l}{\partial \theta}(y_t; a\beta \gamma' y_{t-1}), \\
\frac{\partial \log l}{\partial a}(y_t|y_{t-1};A) = \gamma' y_{t-1} \beta' \frac{\partial \log l}{\partial \theta}(y_t; a\beta \gamma' y_{t-1}).\n\end{cases} (4.6)
$$

The asymptotic properties, especially the asymptotic distribution of the ML estimators of  $\alpha, \beta, \gamma$ , depend on the location of the true vectors  $\beta_0, \gamma_0$ . These properties are straightforward under the strict positivity assumption below.

**Strict Positivity Assumption A.2 iii):** The entries of  $\beta_0$  and  $\gamma_0$  are strictly positive.

Under this strict positivity assumption, the non-negativity-constrained ML estimators have asymptotically strictly positive entries, and the unconstrained and nonnegativity constrained estimators are asymptotically equivalent. However, the ML estimator has to account for the linear constraint of unit sum. This estimator without the non-negativity restrictions is defined as :

$$
(\hat{a}, \hat{\beta}, \hat{\gamma}) = \arg \max_{a, \beta, \gamma} \sum_{t=1}^{T} \log l(y_t; a\beta \gamma' y_{t-1}), \ \text{s.t} : \beta' e = \gamma' e = 1.
$$

The first-order conditions for the Lagrange multipliers associated with the linear restrictions and denoted by  $\lambda, \mu$ , are :

$$
\sum_{t=1}^{T} [\alpha \gamma' y_{t-1} \frac{\partial \log l}{\partial \theta} (y_t; \alpha \beta \gamma' y_{t-1})] - \lambda = 0, \sum_{t=1}^{T} [\alpha y_{t-1} \beta' \frac{\partial \log l}{\partial \theta} (y_t; \alpha \beta \gamma' y_{t-1})] - \mu = 0,
$$
  

$$
\sum_{t=1}^{T} [\gamma' y_{t-1} \beta' \frac{\partial \log l}{\partial \theta} (y_t; \alpha \beta \gamma' y_{t-1})] = 0, \quad \beta' e = \gamma' e = 1.
$$

The FOC need to be solved in  $a, \beta, \gamma, \lambda, \mu$ .

Asymptotically, when T tends to infinity, we get consistent and asymptotically normal ML estimators of B and C. Their asymptotic variance-covariance matrix has a standard form [see Gourieroux and Monfort (1995), Section 10.3] and its estimate can be computed by standard software.

**Remark 4:** As pointed out in Section 3.2.1,  $\beta_1$  (resp.  $\gamma_1$ ) can be interpreted as an eigenvector of A (resp. A'). It is easy to check that  $\beta_1$  (resp.  $\gamma_1$ ) is an eigenvector of AA' (resp. A'A). This implies that  $A = \beta_1 \gamma_1'$  is a singular value decomposition (SVD) of matrix A, for which statistical inference is available mainly in a Gaussian framework [Anderson (1963), Tipping and Bishop (1999)]. However, for  $Rk(A_0) = Rk_+(A_0) = 1$ , the SVD estimation method is not relevant, as its standard asymptotic properties do not account neither for the non-negativity of the data nor the non-negativity of matrix  $A_0$ . Moreover, the SVD interpretation of the NMF is no longer valid for  $Rk_{+}(A_0) \geq 2$ . Indeed, matrix  $AA'$  (resp.  $A'A$ ) is also non-negative, and, except the Perron, Froebenius eigenvector  $\beta_1$  of  $AA'$ , all other eigenvectors  $\beta_2, \beta_3, \dots$  must have at least one negative, or non-real component.

# 4.3 Nonnegative rank  $K_0 = Rk_+(A_0) \geq 2$ .

As mentioned in Section 3.2.3, it is sufficient to estimate one of the admissible NMF's, i.e. an origin in the identified set, to deduce the sharp identified set. The problem with applying the AML method is that the convergence of the AML estimator to the identifed set does not imply its pointwise convergence to a given NMF, which could be used as an origin. Moreover, since the origin is used to derive the asymptotic distribution of the estimator of the identified, it has to be in the interior of the identified set. This section shows how an IML algorithm can solve these issues, allowing us to derive the asymptotic distribution of the estimated set of admissible NMF's. For expository purpose, we provide in the text the relevant assumptions. The additional assumptions needed for asymptotic analysis are given in online Appendix 2.

#### 4.3.1 Consistency of the alternating ML estimator

In this section we consider the consistency of AML approximation of the identified set when T tends to infinity and the number of iterations  $p_T$  in the AML algorithm depends on  $T$  in a suitable manner. Let us consider the dynamic model :

$$
l(y_t|y_{t-1};A) = l(y_t; a \sum_{k=1}^{K} \pi_k \beta_k^* \gamma_k^{*'} y_{t-1}),
$$
\n(4.7)

with  $\beta_k^{*'}$  $k^{*'}e = \gamma_k^{*'}$  $k \neq k' = 1, k = 1, \ldots, K, \pi' e = 1$  and the identified set:

$$
\mathcal{A}_0 = \{a, \pi_k, \beta_k^*, \gamma_k^*, k = 1, \dots, K, \text{ such that } a \sum_{k=1}^K \pi_k \beta_k^* \gamma_k^{*'} = A_0 \}.
$$

where  $A_0$  is the true value of A.

Section 2.2 shows that the underlying parametric families  $l(y; \theta), \theta \geq 0$ , are often constructed from the products of Poisson, or exponential distributions. Therefore, they satisfy the following assumption :

**Assumption A.4**: The underlying log-likelihood log  $l(y; \theta)$  is concave in  $\theta, \theta \ge 0$ .

Under Assumption A.4, if  $l(y_t|y_{t-1}, A) = l(y_t, Ay_{t-1})$  and the components of  $y_t$  are non-negative, each step of the AML algorithm outlined in Section 4.1.2 leads to a unique solution in  $(B, C)$ , because the objective function is log-concave in B (resp. C) for a given C (resp.  $B$ ), and then also under the alternative parametrization  $(a, \pi_k, \beta_k^*, \gamma_k^*, k = 1, \dots, K)$ . Therefore, the AML estimator is a function of the underlying (normalized) log-likelihood  $\frac{1}{\pi}$  $\frac{1}{T}L_T(A) = \frac{1}{T}$  $\sum$ T  $t=1$  $log l(y_t; Ay_{t-1})$ , and of the initial values used in the algorithm (and of the number of iterations  $p$ ). Let the set of parameters be denoted by  $\alpha = (\pi_k, \beta_k^*, \gamma_k^*, k = 1, ..., K)$  and the selected initial value

by  $\alpha^o$ . The AML estimator at iteration p can be written as :

$$
\hat{\alpha}_T(\alpha^o, p) \equiv \delta(\frac{1}{T}L_T(.); \alpha^0, p),\tag{4.8}
$$

where  $\delta$  is a deterministic function. Then, for large T, the AML estimator will converge asymptotically to the value :

$$
\alpha_{\infty}(\alpha^o, p) = \delta[E_0 \log l(Y_t; A_0 Y_{t-1}); \alpha^0, p],\tag{4.9}
$$

that belongs in the set  $\mathcal{A}_0$  if p is large. Note that this limiting value can depend on the starting value  $\alpha^0$  of the algorithm. More precisely, under Assumptions A.1-A.4 and the additional regularity conditions a.1 given in online Appendix 2, we have the following proposition as a direct consequence of the numerical consistency of the AML approximation (for  $p \to \infty$  and T fixed) and of the uniform convergence in Proposition 3 :

**Proposition 4:** For large T, there exist a function  $c(.)$  and a number of iterations  $p_T$  such that, for any  $p \geq p_T : \mathcal{D}[\alpha_T(\alpha^o, p), \mathcal{A}_o] < c(\alpha^o)/T$ , where  $\mathcal D$  measures the distance between  $\alpha_T(\alpha^o, p)$  and the set  $\mathcal{A}_o$ .

In practice, we can apply the AML algorithm with a given starting value  $\alpha^o$  and a number of iterations  $p$ . The number  $p$  needs to be set sufficiently large for Proposition 4 to be satisfied. Then, the asymptotic bias of the alternating ML estimator will be sufficiently small to become negligible in the asymptotic distribution of the IML estimator derived later in the text. Two optimizations are performed at each step of the algorithm, with respect to  $B$  and  $C$ , respectively. A Newton-Raphson type of algorithm can be used in each of these optimizations. Under the log-concavity assumption A.4, the Newton-Ralphson algorithm is a special case of the steepest ascent algorithm, which increases the objective function at each step. As the increase of the objective function at each step of the algorithm is ensured, the AML with a fixed number of iterations in each intermediate optimization will also increase the objective function, which is sufficient for the numerical consistency of the AML algorithm under Assumption A.4.

It follows from Proposition 1 that all other elements of  $\mathcal{A}_0$  are functions of  $\alpha_\infty(\alpha^o; p)$ 

and of the elements of a matrix Q with unitary diagonal elements, such that :

$$
B[\alpha_{\infty}(\alpha^o, p)]Q \ge 0, \quad C[\alpha_{\infty}(\alpha^o, p)][Q']^{-1} \ge 0.
$$

This defines a set  $\mathcal{Q}[\alpha_{\infty}(\alpha^o, p)]$  of admissible values of transformation Q.

Equivalently, we have a parametric representation of the identified set  $\mathcal{A}_0$  of NMF's :

$$
\mathcal{A}_0 = \{ \alpha : \alpha = \zeta[\alpha_\infty(\alpha^\circ, p), Q], Q \in \mathcal{Q}[\alpha_\infty(\alpha^\circ, p)] \},\tag{4.10}
$$

where  $\zeta$  is a known function. Then, the set  $\mathcal{A}_0$  is consistently estimated as:

$$
\hat{\mathcal{A}}_T = \{ \alpha : \alpha = \zeta[\hat{\alpha}_T(\alpha^o, p), Q], Q \in \mathcal{Q}[\hat{\alpha}_T(\alpha^o, p)] \equiv \hat{Q}_T(\alpha^o, p) \}. \tag{4.11}
$$

The estimation method presented above will allow us to approximate the identified set for a sufficiently large  $p$ .

Proposition 4 implies that the above AML (set) estimator is a set maximizer of the log-likelihood function, which is flat at the maximum. This AML estimator could be used to derive a Monte-Carlo confidence set for the identified set based on a quasilikelihood ratio with the quantiles computed by simulations [see, Chen, Christensen] and Tamer (2018), Remark 1, p. 1972]. In our framework, the chosen quasi-prior distribution would need to ensure the existence of the log-likelihood function  $L_T(A)$ . For example, it would have to satisfy the non-negativity restrictions on matrix A in a Poisson autoregressive model. Moreover, the generic form of the estimated confidence set, i.e.  $\{A : L_T(A) \geq \xi\}$ , where  $L_T$  denotes the log-likelihood and  $\xi$  is the estimated quantile, is hard to determine unless the appropriate parametrization given in Proposition 1 is used.

Because of the lack of identification, we do not have the numerical stability of  $\hat{\alpha}_T(\alpha_0, p)$ , for large p. Therefore, we cannot expect to prove the asymptotic normality of this AML estimator. In order to stabilize the algorithm, we include an additional optimization step.

## 4.3.2 How to introduce an identification restriction by Identifying Maximum Likelihood

An identification issue is commonly solved either by introducing implicit identification restrictions, or by reparametrizing the model and dividing the parameters into the set of identifiable and non-identifiable parameters (this is the "global" reduced form reparametrization considered in Shi and Shum (2015) and Chen et al. (2018), Section 5.1.1). These approaches are not suitable in our framework, where the parametrization of the identified set depends on a selected origin in  $A_0$ . Hence, we introduce indirectly  $K(K-1)$  identification restrictions.

Let us consider an alternative parametrization method. For a given origin, the identifiable set is parametrized by  $vec^*Q = q$ , where  $vec^*Q$  denote the stacked elements of Q except for the diagonal elements equal to 1. Let  $\tilde{\alpha}$  denote the solution of the AML algorithm at step p.  $\tilde{\alpha}$  is not necessarily an element in the interior of the identified set. Moreover, it depends on the starting value.

Let us now determine an interior origin  $\alpha_0^*$ . We define the extended parameter vector  $(q, \tilde{\alpha}) \equiv \alpha$ . Using this notation, the NMF  $\tilde{\alpha}$  is now represented as  $(0, \tilde{\alpha})$ . We define the interior origin of the identified set as the optimizer of an additional criterion with respect to the additional parameter q. Several criteria  $\tilde{g}(q, \tilde{\alpha}) = g(\alpha)$ , where  $\alpha \in \mathcal{A}_0$ , arise naturally :

i) The concentration of a discrete probability distribution is usually measured by  $\sum$ K  $k=1$  $(\pi_k \log \pi_k)$ . This quantity is negative, attains its minimum for  $\pi_k = 1/K, \forall k$ , that is a uniform distribution, and increases to zero with the concentration of the distribution. Therefore, given  $\tilde{\alpha}$ , the factorization provides us the least concentrated latent heterogeneity defined by :

$$
\alpha_0^* = \arg \max_{\alpha \in \mathcal{A}_0} - \sum_{k=1}^K \pi_k(q, \tilde{\alpha}) \log \pi_k(q, \tilde{\alpha}),
$$

where the minimization is in q for a given  $\tilde{\alpha}$ . This means that the new selected origin is  $\alpha_0^* = (q(\tilde{\alpha}), \tilde{\alpha})$ , where:

$$
q^*(\tilde{\alpha}) = Argmax_{q \in Q(\tilde{\alpha})} - \sum_{k=1}^K \pi_k(q, \tilde{\alpha}) \log \pi_k(q, \tilde{\alpha}),
$$

where the domain of q can depend on  $\tilde{\alpha}$ .

ii) The volume/collinearity criterion is:  $g(\alpha) = \det(B^{*'}B^*)$ , where  $B^* = (\beta_1^*, \dots, \beta_K^*)$ , or  $g(\alpha) = \det(C^{*'}C^*)$ , where  $C^* = (\gamma_1^*, \ldots, \gamma_K^*)$ . As above  $B^*, C^*$  depend on  $(q, \tilde{\alpha})$ , but we do not account for it to keep the notation simple. The above criterion measures

the volume of the parallelepiped generated by the columns of  $B^*$  (resp.  $C^*$ ). The larger this volume, the less "collinear" the columns of  $B^*$  (resp.  $C^*$ )<sup>[14](#page-0-0)</sup>. Alternatively, the collinearity measure could be applied to the matrix  $B$  itself. Let  $\tilde{B}$  denote the corresponding element of  $\tilde{\alpha}$ . Then we have:

$$
B(q,\tilde{\alpha}) = \tilde{B}Q, \text{ and } det[B(q,\tilde{\alpha})'B(q,\tilde{\alpha})] = det[Q'\tilde{B}'\tilde{B}Q] = (det(Q))^2 det(\tilde{B}'\tilde{B}).
$$

The maximization of the collinearity measure is equivalent to the maximization of  $(det(Q))^2$  on the domain  $Q(\tilde{\alpha})$ . This objective function is independent of  $\tilde{\alpha}$ , but the constraints on q depend on  $\tilde{\alpha}$ .

The patterns of these criteria and their combinations are illustrated and discussed in Online Appendix 3 for the example of Section 6. When they are used, the optimum of function g can still be reached on the boundary of the domain of values of  $q_{1,2}, q_{2,1}$ rather than in its interior.

iii) The repulsion criterion. It follows from Proposition 1 that the points in the interior of the identified set are such that the elements  $\beta_{i,j}^*$  and  $\gamma_{i,j}^*$  are strictly positive. Therefore, a point in the interior can be reached by combining the previous criteria with a third one such that  $\sum_i \sum_j (ln(\beta^*_{i,j}) + ln(\gamma^*_{i,j}))$  that creates a repulsion effect on the boundary of the domain.

To approximate the interior origin  $\alpha_0^*$ , an additional optimization step needs to be included in the AML algorithm<sup>[15](#page-0-0)</sup>. Let the recursive system in this algorithm be denoted by:  $\alpha^{(p+1)} = H(\alpha^{(p)})$ , where H depends on the observations. Then, the identifying maximum likelihood (IML) algorithm is the following :

step 1: Select an initial value  $\alpha^{(0)}$ . The parameter value needs to be in the parameter set. In particular, the associated starting values  $B^{(0)}$ ,  $C^{(0)}$  have to be nonnegative and of a given rank  $K$ . This rank condition on the starting value is important to ensure non-degenerate behavior of the algorithm at the next steps.

step p: At step p, a value  $\alpha^{(p)}$  is available.

i) Apply the AML algorithm to get a value  $\tilde{\alpha}^{(p+1)} = H(\alpha^{(p)})$ .

This value is considered as an approximation of a point in the identifiable set.

<sup>&</sup>lt;sup>14</sup>These criteria are the analogues of SVD identification restrictions:  $B'B = C'C = Id$ , where all factorial directions are orthonormal.

<sup>&</sup>lt;sup>15</sup>This additional step is the analogue of step 2 in the estimation approach introduced in Davezies et al. (2022), Section 3.1., where it is applied to an approximate identified set instead of the identified set itself.

It can be used to parametrize the set  $\hat{\mathcal{A}}^{(p)}$  by q.

ii) Perform the optimisation of the additional criterion to get :

$$
q^{(p+1)} = \text{Opt}_{q \in Q^{(p+1)}} \tilde{g}(q; \tilde{\alpha}^{(p+1)}),
$$

where  $Q^{(p+1)}$  is the domain defined by the inequality restrictions with  $\tilde{\alpha}^{(p+1)}$ . Then, the solution is a function of  $\tilde{\alpha}^{(p+1)}$ , that is :  $q^{(p+1)} = q(\tilde{\alpha}^{(p+1)})$ , say, where function  $q(.)$  does not depend on the observations.

iii) Find  $\alpha^{(p+1)}$  by transforming  $\tilde{\alpha}^{(p+1)}$  with the linear transformation  $Q^{(p+1)}$  associated with  $q^{(p+1)}$  to get  $\pi_k^{(p+1)}$  $\binom{(p+1)}{k}$ ,  $k = 1, ..., K$  arranged in a decreasing order<sup>[16](#page-0-0)</sup>. More precisely:

compute  $Q^{(p+1)}$  such that  $q^{(p+1)} = vec^*Q^{(p+1)}$ ; compute  $\tilde{B}^{(p+1)}$ ,  $\tilde{C}^{(p+1)}$  from  $\tilde{\alpha}^{(p+1)}$ ; compute  $B^{(p+1)} = \tilde{B}^{(p+1)} Q^{(p+1)}$ ,  $C^{(p+1)} = \tilde{C}^{(p+1)} [Q^{(p+1)'}]^{-1}$ ; compute  $\alpha^{(p+1)}$  from  $B^{(p+1)}$ ,  $C^{(p+1)}$ , etc.

The main difference between the AML and IML algorithms concerns the consistency when  $T, p_T$  tend to infinity according to Proposition 4. The AML converges to the set  $A_0$ , but the convergence is not pointwise. The IML converges to a given interior origin  $\alpha_0^*$  that allows us to perform a Taylor expansion of the first-order conditions in order to derive the asymptotic normality, as in Shi and Shum (2015).

The IML method requires the availability of algorithms for the optimization of nonlinear functions under a large number of nonlinear inequality restrictions. The recent developments in Active Set Sequential Quadratic Programming (SQP) have largely solved this problem [see e.g. Gill et al. (2002) and Liu (2005) for a proof of numerical convergence].

Remark 5: The additional intermediate optimization in the IML algorithm does not necessarily have to be introduced starting from the first iteration  $p = 1$ . It can be introduced later, when  $p$  is sufficiently large to get a value close to the identified set by Proposition 4. In this respect, the IML is used to stabilize pointwise the values obtained from the AML algorithm.

<sup>16</sup>To solve the up-to-permutation idenification issue.

#### 4.3.3 Asymptotic distributions

The standard asymptotic arguments can be used to derive the asymptotic normality of the IML estimator adjusted to reach  $\alpha_0^*$ , since  $\alpha_0^*$  is in the interior of  $\mathcal{A}_0$  and if the associated  $\pi_{0k}^*, k = 1, ..., K$  are all distinct<sup>[17](#page-0-0)</sup>. To clarify the role of the intermediate optimization in the IML algorithm, let us first consider the standard information matrices. Two information matrices appear naturally<sup>[18](#page-0-0)</sup>:

i) The information based on unconstrained  $A$ , is  $E_0$  $\lceil$  $-\frac{\partial^2 \log l(y_t|y_{t-1};A)}{2}$  $\partial$  vec  $A \partial$  vec  $A'$ 1 . This matrix is invertible by the assumption of identifiable A, but can be of a high dimension in practice.

ii) The information matrix corresponding to parameter  $\alpha$ :  $J_0 = E_0 \left[-\frac{\partial^2 \log l(y_t|y_{t-1}; \alpha)}{\partial x \partial y}\right]$  $\frac{\partial \alpha \partial \alpha'}{\partial \alpha'}$ is constrained by the unit mass restrictions on  $\pi, \beta_k^*, \gamma_k^*, k = 1, \ldots, K$ . This matrix has a smaller dimension, but is not of full rank because of the lack of identification.

The IML algorithm introduced above has extended the constrained maximum likelihood approach by adding asymptotically the identification restrictions corresponding to the first-order conditions of the optimization of  $\tilde{g}(q,\alpha)$  with respect to q, that are :

$$
\frac{\partial \tilde{g}(q,\alpha)}{\partial q} = 0 \Rightarrow q = q(\alpha),\tag{4.12}
$$

of a number equal to the degree of underidentification.

These limiting conditions have been replaced by  $\frac{\partial \tilde{g}}{\partial \theta}$  $\frac{\partial g}{\partial q}[q(\hat{\alpha}_T), \hat{\alpha}_T] = 0$  in the IML algorithm and can be expanded in a neighbourhood of  $[q(\alpha_0^*) = 0, \alpha_0^*]$ . We get :

$$
\left[\frac{\partial^2 \tilde{g}}{\partial q \partial q'}[0, \alpha_0^*]\frac{dq}{d\alpha}(\alpha_0^*) + \frac{\partial^2 \tilde{g}}{\partial q \partial \alpha'}(0, \alpha_0^*)\right] \sqrt{T}(\hat{\alpha}_T - \alpha_0^*) \simeq 0,
$$

that are the additional asymptotic linear restrictions  $D_2'$  $\overline{T}(\hat{\alpha}_T - \alpha_0^*) = 0$  on  $\sqrt{T}(\hat{\alpha}_T \alpha_0^*$ ), say, to be taken into account for the computation of the asymptotic variance of  $\hat{\alpha}_T$ . More precisely,

<sup>&</sup>lt;sup>17</sup>This additional condition is analogous to the condition of distinct eigenvalues in the joint spectral decomposition of  $AA'$  and  $A'A$  in the SVD.

<sup>&</sup>lt;sup>18</sup>For expository purpose, we keep the same notation  $l$  for the conditional likelihood as a function of A, or a function of  $\alpha$ .

**Proposition 5**: If  $\alpha_0^*$  is in the interior of  $\mathcal{A}_0$ , if the associated  $\pi_{0k}^*$ ,  $k = 1, ..., K$ , are all distinct, if assumptions A.1-A.4 and the additional regularity assumptions a.1-a.2 are satisfied, then,  $\sqrt{T}(\hat{\alpha}_T - \alpha_0^*) \stackrel{d}{\longrightarrow} N(0, J_0^{11} J_0 J_0^{11})$ , where  $J_0^{11}$  is the North-West block in the block decomposition of the inverse of matrix  $\begin{pmatrix} J_0 & D \\ D_0 & 0 \end{pmatrix}$  $D'$  0  $\setminus$  $, D =$  $(D_1, D_2), D_1$  defining the  $2K+1$  unit mass restrictions on  $\pi_k, \beta_k^*, \gamma_k^*, k = 1, ..., K$ , and  $D_2$  defining the (asymptotic) linearized restrictions corresponding to the intermediate optimizations in the IML algorithm.

Proof: See Appendix 1.1.

In the expression of the asymptotic variance-covariance matrix, the three main components are the information related with the unconstrained ML of  $\alpha$ , the unit mass restrictions  $D_1$  and the restrictions  $D_2$  due to additional intermediate optimization, respectively.

Remark 6: The sequence of optimizations cannot be replaced by a penalty term in the objective function as suggested by the machine learning literature either to ensure the (numerical) convergence of the AML algorithm [see e.g. Kim and Park (2008), Schachtner et al. (2011), Hastie et al. (2015)], or to circumvent the curse of dimensionality [De Paula et al. (2023), eq. (10)]. In our framework, this Penalty Function Approach (PFA) would lead to an objective function of the form  $\log l_T(y; \alpha) + \lambda_T g(\alpha)$ , where the tuning parameter  $\lambda_T$  would be a function of T, appropriately chosen to ensure the numerical convergence. It is easy to see why this approach will not provide an estimator converging to an interior origin in the identified set. Indeed, the asymptotic first-order conditions would become  $\frac{\partial \log l_T(y,\alpha)}{\partial \alpha}$  $rac{\partial q}{\partial \alpha} + \lambda_T$  $\partial g(\alpha)$  $rac{g(\alpha)}{\partial \alpha}$ . These FOCs are not aligned with the direction allowing us to remain in the set, since  $\frac{\partial g(\alpha)}{\partial \alpha}$  $rac{\partial(\alpha)}{\partial \alpha} =$  $d\tilde{g}[q(\alpha),\alpha]$  $d\alpha$ differs from  $\frac{\partial \tilde{g}(q,\alpha)}{\partial}$  $\frac{\partial q}{\partial q}$  [see also Ariaz et al. (2018), Section 5, for a critique of PFA in a SVAR model with partial identification].

The asymptotic Gaussian uncertainty on the IML estimator of  $\alpha_0^*$  determines all the uncertainties on the IML estimator of the set  $\mathcal{A}_0$  of NMF's. More precisely, any other element of  $\mathcal{A}_0$  can be written as a deterministic function of  $\alpha_0$ :  $\xi(q, \alpha_0^*)$ , with  $q \in Q(\alpha_0^*)$ . Then, that element can be estimated by  $\xi(q, \hat{\alpha}_T)$ , which is a given function of  $\alpha_T$ . Therefore, it inherits the asymptotic properties of  $\hat{\alpha}_T$ : it is consistent of  $\xi(q, \alpha_0^*)$ , asymptotically normal, and its asymptotic variance-covariance matrix is obtained from the Slutsky formula (i.e. the  $\delta$ -method), whenever q is not on the boundary of  $Q(\alpha_0^*)$ . If q is on the boundary, its asymptotic distribution will become a truncated normal, and can be easily found by simulation.

**Remark 7:** The IML approach can be used to derive the lower and upper bounds on partially identified scalar parameters, and to obtain the measures of uncertainty on these bounds. The examples are the minimum and maximum values of functions  $\sum$ K  $k=1$  $(\pi_k \log \pi_k)$ , det $(\tilde{B}'\tilde{B})$ , det $(\tilde{C}'\tilde{C})$ . This procedure is analogous to determining the confidence intervals for average marginal effects in a fixed effects panel logit model [Liu et al. (2023), Davezies et al. (2022)]. Likely, there exists an interval of admissible values of the above uncertainty measures. Such an interval is easy to obtain for the measure of concentration when  $K = 2$ . However, the lower and upper bounds will be reached on the boundaries of the identified set for  $q_{12}, q_{21}$  when  $K = 2$ , for example. Then, the joint asymptotic distribution of these bounds cannot be Gaussian due to the effect of infimum in Proposition  $2^{19}$  $2^{19}$  $2^{19}$ . Note that the joint asymptotic distribution of these two bounds is easily derived by simulations and, by construction, we cannot have bound reversal in the estimation.

**Remark 8:** The IML approach can also be used to derive a confidence set for  $\mathcal{A}_0$  with the correct asymptotic level, following Shi and Shum (2015). However this asymptotic confidence set is difficult to represent graphically due to the high dimension of  $\mathcal{A}_0$ .

## 4.3.4 Asymptotic distribution of  $\hat{A}_T$

The IML approach helps us find the estimates and confidence intervals of identifiable parameters, such as the elements  $a_{ij}$  of matrix A. In practice, we may encounter the curse of dimensionality that complicates the unconstrained estimation of A. Moreover, the estimator has to be applied under the constraint of a given non-negative rank :  $Rk_{+}(A) = K$ , and the rank-constrained confidence intervals are likely narrower than the unconstrained ones. They can be derived from  $\hat{\alpha}_T$  by simulations, given that :

 $19$ In this respect, Assumption A.1 i) of normality of the upper and lower bounds in Imbens et al. (2004), Stoye (2009) is not satisfied in our framework.

$$
a_{ij} \simeq \hat{a} \sum_{k=1}^K \hat{\pi}_k \hat{\beta}_{ik}^* \hat{\gamma}_{jk}^* = \hat{a}_{ij}.
$$

Asymptotically,  $\hat{a}_{ij}$  will converge to the true value  $a_{ij,0}$  that does not depend on the choice of interior origin  $\alpha_0^*$  . Similarly, its asymptotic variance-covariance matrix is also independent of this choice, i.e. of the selected function  $\tilde{q}$  and of the starting values of the IML algorithm. The additional constraints  $D_2$  are introduced only for solving the identification issue.

Proposition 6 The asymptotic distribution of the IML estimator of matrix A does not depend on the interior origin  $\alpha_0^*$  in the identified set, i.e. on the choice of the additional optimization criterion.

Proof: See Appendix 1.2.

We have introduced an IML estimator of  $A_0$  under the non-negative rank restriction  $Rk_{+}(A_0) = K_0$ , derived its asymptotic Gaussian behavior<sup>[20](#page-0-0)</sup> and found the expression of the associated efficiency bound on this identifiable matrix parameter. The asymptotic variance-covariance matrix of  $\hat{A}_T$  is obtained by applying the Slutsky formula based on the first-order expansions of  $a \sum_{k=1}^{K} \pi_k \beta_k^* \gamma_k^*$  in a neighbourhood of  $\alpha_0^* = (a_0^*, \pi_{0k}^*, \beta_{0k}^*, \gamma_{0k}^*, k = 1, ..., K)$ . Then, the asymptotic confidence intervals for the identifiable parameter functions of A also share the asymptotic efficiency properties.

The asymptotic expansion and normality of the IML estimator of A could also be used to derive the asymptotic distribution of the log-likelihood ratio test statistic corresponding to the test of hypotheses:  $H_0: Rk_+A = K - 1$  against  $H_0: Rk_+A = K$ that is a chi-square distribution with an appropriate degree of freedom. Then, we would deduce a sequence of test procedures to estimate  $Rk_{+}(A)$ . Due to the curse of dimensionality we cannot follow a "general to specific" approach, starting from the highest possible rank  $K = inf(n, m)$  and testing it against the alternative of a lower value of K. Instead, we have to start from the smallest value of  $K = 0$  and test it against higher values of  $K = 1, 2, ...$ 

<sup>&</sup>lt;sup>20</sup>This asymptotic Gaussian distribution is degenerate because of the reduced rank.

## 5 Extension to Nonparametric Identification

### 5.1 The Identified Set

The proof of Proposition 1 is valid in a nonparametric framework. Let us consider a pair of real variables  $(X, Y)$  with a continuous joint distribution on a product of intervals and assume a positive pdf  $f(x; y)$  on that interval. We can be interested in a decomposition :

$$
f(x,y) = \sum_{k=1}^{K} \pi_k \beta_k(x) \gamma_k(y), \qquad (5.1)
$$

where  $(\pi_k)$  defines the latent heterogeneity distribution, and  $\beta_k(x)$ ,  $\gamma_k(y)$  are the pdf's of  $x$  and  $y$ , respectively. This is a mixture model in which each joint distribution  $\beta_k(x), \gamma_k(y)$  in the mixture satisfies the independence condition [see Compiani and Kitamura (2016) for mixtures in Econometrics $]^{21}$  $]^{21}$  $]^{21}$ .

The proof remains valid if functions  $\beta_1, \ldots, \beta_K$  (resp.  $\gamma_1, \ldots, \gamma_K$ ) are a.s. linearly independent. Since the model is already normalized, Proposition 1 becomes :

**Proposition 1'**: Let us consider a true factorization of the joint pdf :

$$
f_0(x, y) = \sum_{k=1}^{K} b_{0k}(x) \gamma_{0k}(y) = B_0(x) C_0(y),
$$

where the functions  $b_{0k}$  are positive densities (not necessarily with unitary mass) and  $\gamma_{0k}$  are probability densities. Then, the other observationally equivalent decompositions are  $B(x) = B_0(x)Q$ ,  $C(y) = C_0(y)(Q')^{-1}$ , where the matrix Q is invertible, with unitary diagonal elements and such that  $:B_0(x)Q \geq 0$ , a.s. in x,  $C_0(y)(Q')^{-1} \geq$ 0 a.s. in y.

The decomposition (5.1) of the joint p.d.f. has alternative interpretations. For instance, the conditional p.d.f. can be written as:

$$
f(y|x) = \frac{f(x,y)}{f(x)} = \sum_{k=1}^{K} \frac{b_{0k}(x)}{f(x)} \gamma_{0k}(y) \equiv \sum_{k=1}^{K} b_{0k}^{*}(x) \gamma_{0k}(y).
$$

<sup>&</sup>lt;sup>21</sup>This condition can be written for more than two variables, so that the mixture becomes identifiable and easier to analyse [Hall and Zhou (2003), Kasahara and Shimotsu (2009)]. A suitable notion of nonnegative rank has not yet been introduced in the literature for 3- or 4-entry tables.

Therefore, it is equivalent to impose a reduced rank condition on the joint distribution, or on the conditional distributions. This problem is considered in Henry, Kitamura and Salanie  $(2014)$ , where the pair of variables is denoted by Y, W, the representation is written conditionally on a third variable  $X$ , and the independence condition on the elements of the mixtures is called the exclusion restriction [see also Compiani and Kitamura (2016)].

### 5.2 Nonparametric Identifying Maximum Likelihood

The estimation approach outlined in Section 4 can be extended to the functional parameter framework, with functional parameters  $b_k, k = 1, \ldots, K$ , scalar parameters  $\pi_k, k = 1, \ldots, K$ , and additional parameters in  $vec^*Q$ . Let us assume i.i.d. observations on  $(X_i, Y_i)$ ,  $i = 1, ..., n$ . Then, a kernel based AML can be applied with the objective function :

$$
\sum_{i=1}^{n} \{ K(\frac{x_i - x}{h_n}) K(\frac{y_i - y}{h_n}) \log \sum_{k=1}^{K} [b_k(x)\gamma_k(y)] \},\
$$

where K(.) denotes a kernel,  $h_n$  the bandwidth and  $\int \gamma_k(y)dy = 1, k = 1, ..., K$ . The maximisation with respect to  $b_k(x), \gamma_k(y), k = 1, ..., K$  provides the functional estimators of  $b_k(.)$ ,  $\gamma_k(.)$ ,  $k = 1, ..., K$ , and then functional estimators of  $\pi_k, \beta_k(.)$ ,  $\gamma_k(.)$ .

Next, to solve the identification issue, an additional optimisation has to be performed to fix a mixture in the identified set. Criteria equivalent to the criteria introduced in Section 4.3.2 can be used, either the concentration criterion  $\sum$ K  $k=1$  $(\pi_k \log \pi_k),$ or the volume/collinearity criterion det Γ, with  $\Gamma = \int \gamma(y) \gamma'(y) dy$ .

The analysis of the asymptotic properties of this functional estimation approach are out of the scope of this paper and are left for future research<sup>[22](#page-0-0)</sup>.

 $^{22}$ In this framework the benchmark is a functional parameter, a case that is not included in Shi and Shum (2015).

# 6 Illustration

We consider in this section observations  $y_t, t = 1, \ldots, T$ , on a process satisfying a multivariate autoregressive conditional Poisson (ACP) model:

$$
y_{it}|y_{t-1} \sim \mathcal{P}(a_i y_{t-1} + \mu_i), \quad a_i \ge 0, \ \mu_i \ge 0, \ i = 1, \dots, n,
$$
 (6.1)

where  $y_{1t}, \ldots, y_{nt}$  are independent given  $y_{t-1}$ . The process is such that  $E(Y_t|Y_{t-1}) =$  $AY_{t-1} + \mu$  and ergodic if the eigenvalues of A are of modulus strictly less than 1.

## 6.1 The true dynamics

We consider a process of dimension four  $n = 4$ , and the true matrix  $A_0$  of non-negative rank equal to 2. This matrix is given by :

$$
A_0 = \frac{1}{96} \begin{pmatrix} 1 \\ 1 \\ 1 \end{pmatrix} (1, 1, 2, 2) + \frac{1}{96} \begin{pmatrix} 1 \\ 2 \\ 1 \end{pmatrix} (1, 1, 1, 1)
$$
  
= 
$$
\frac{1}{96} \begin{pmatrix} 2 & 2 & 3 & 3 \\ 3 & 3 & 4 & 4 \\ 2 & 2 & 3 & 3 \\ 3 & 3 & 4 & 4 \end{pmatrix} = \begin{pmatrix} 0.0208 & 0.0208 & 0.0312 & 0.0312 \\ 0.0312 & 0.0312 & 0.0416 & 0.0416 \\ 0.0208 & 0.0208 & 0.0312 & 0.0312 \end{pmatrix}.
$$
 (6.2)

and the true intercept is  $\mu_0 = (2, 2, 2, 2)'$ . The matrix  $A_0$  can be rewritten to highlight its interpretations in terms of probability distributions (see eq. 2.13). We have :

$$
A_0 = a_0 [\pi_0 \beta_{01}^* \gamma_{01}^{*'} + (1 - \pi_0) \beta_{02}^* \gamma_{02}^{*'}], \tag{6.3}
$$

where  $a_0 = 1/2, \pi_0 = 1/2$ ,

$$
\beta_{01}^* = \frac{1}{4} \begin{pmatrix} 1 \\ 1 \\ 1 \\ 1 \end{pmatrix}, \gamma_{01}^* = \frac{1}{6} \begin{pmatrix} 1 \\ 1 \\ 2 \\ 2 \end{pmatrix}, \beta_{02}^* = \frac{1}{6} \begin{pmatrix} 1 \\ 2 \\ 1 \\ 2 \end{pmatrix}, \gamma_{02}^* = \frac{1}{4} \begin{pmatrix} 1 \\ 1 \\ 1 \\ 1 \end{pmatrix}.
$$
 (6.4)

By Proposition 2, we get :

$$
-1 \le q_{12} \le 1, \ \ -1/2 \le q_{21} \le 1/2. \tag{6.5}
$$

Therefore the identified set for the matrix :

$$
B = B_0 Q = \frac{1}{96} \begin{pmatrix} 1 & 1 \\ 1 & 2 \\ 1 & 1 \\ 1 & 2 \end{pmatrix} \begin{pmatrix} 1 & q_{12} \\ q_{21} & 1 \end{pmatrix},
$$

is convex and generated by the four extreme matrices corresponding to :

$$
(q_{12} = -1, q_{21} = -1/2), (q_{12} = -1, q_{21} = 1/2), (q_{12} = 1, q_{21} = -1/2), (q_{12} = 1, q_{21} = 1/2)
$$

For instance the first extreme point is :

$$
B = \frac{1}{96} \begin{pmatrix} 1/2 & 0 \\ 0 & 1 \\ 1/2 & 0 \\ 0 & 1 \end{pmatrix},
$$

and has some zero entries, i.e. satisfies the sparsity property.

We generate simulated trajectories of a four-dimensional count process for  $t =$  $1, \ldots, T = 500$  and display in Figure 1 the first 200 values. The initial values are set equal to  $y_{it} = 3, i = 1, ..., 4$  at  $t = -50$ . The observed values vary between 0 and 10.



Figure 1: Evolution of  $y_t$ 

By construction, the effect of lagged values is captured by the sufficient statistics:  $z_{1,t} = y_{1,t} + y_{2,t}$  and  $z_{2,t} = y_{3,t} + y_{4,t}$  displayed in Figure 2.



Figure 2: Evolution of the Sufficient Statistics

The process  $z_t = (z_{1,t}, z_{2,t})'$  follows a conditional Poisson model with a NMF of nonnegative rank equal to 2, i.e. the dimension of  $z$ . We observe clustering of small and large values in Figures 1 and 2 because of serial dependence of both processes. The true parameter values chosen in this experiment can complicate the inference because matrix  $A_0$  has two non-negative singular values, the first one 37.9 times bigger than the second one, making it difficult to distinguish between the ranks  $Rk(A)=1$  and  $Rk(A)=2$  in practice. Moreover, the uniform distributed heterogeneity in matrix  $A_0$ factorization complicates the identification.

## 6.2 Inference

To find the NMF, we consider the maximum likelihood estimation of matrix A from the simulated data. The log-likelihood computed from the component series  $i =$ 1, ..., 4 at times  $t = 1, ..., T$  and concentrated with respect to  $\mu$  is:

$$
l(A; y) \approx \sum_{i=1}^{4} \sum_{t=2}^{T} \left\{ -(\bar{y}_i + a_i(y_{t-1} - \bar{y})) + y_{i,t} \ln(\bar{y}_i + a_i(y_{t-1} - \bar{y})) \right\}
$$
(6.6)

where  $\bar{y}$  denotes the marginal mean of the series and  $\approx$  stands for the equality up to an additive term independent of parameter A. This concentrated log-likelihood function is maximized by applying the IML algorithm in which the additional maximization criterion combines the concentration measure  $\pi \ln \pi + (1 - \pi) \ln(1 - \pi)$ , the measure of collinearity of the standardized factors  $\beta^*$ , i.e.  $|det(B^{*'}B^*|,$  and the repulsion criterion. The joint use of the three criteria is needed in our experiment. The use of the concentration criterion only in the auxiliary optimization is insufficient to define

an origin in the true identified set. It allows us only to reduce the degree of under identification by one. With the additional collinearity criterion, the solution  $\alpha^*$  is located on the border of the identified set. Then, the repulsion criterion allows us to reach a point  $\alpha_0^*$  in the interior. A discussion of the patterns of these additional objective criteria as a function of  $Q$  is provided in the Online Appendix 3.1 for the non-negative rank  $K = 2$ .

The IML estimated components  $q_{1,2}$  and  $q_{2,1}$  of matrix  $Q$ , based on AML matrices  $\hat{B}_T$  and  $\hat{C}_T$  are reported in Appendix 3.3 and obtained by maximizing the auxiliary objective function with the volume/collinearity and repulsion criteria are  $\hat{q}_{1,2} = 0.1833$ and  $\hat{q}_{2,1} = 0.0872$  for T = 500 and  $\hat{q}_{1,2} = 0.1557$  and  $\hat{q}_{2,1} = -0.0393$  for T = 1000. For the sample size  $T=500$ , we get the IML estimated matrices:

$$
\hat{B}_{500}^{IML} = \hat{B}_{500}^{AMI} * \hat{Q}_{500} = \begin{pmatrix} 0.0291 & 0.1181 \\ 0.0395 & 0.0939 \\ 0.1331 & 0.0279 \\ 0.1182 & 0.0511 \end{pmatrix}, \quad \hat{C}_{500}^{IML} = \hat{C}_{500}^{AMI} * (\hat{Q}_{500}')^{-1} = \begin{pmatrix} 0.0926 & 0.3922 \\ 0.1331 & 0.0279 \\ 0.0490 & 0.0379 & 0.0301 & 0.0232 \\ 0.0405 & 0.0379 & 0.0301 & 0.0232 \\ 0.0233 & 0.0248 & 0.0294 & 0.0628 \\ 0.0310 & 0.0297 & 0.0317 & 0.0580 \end{pmatrix},
$$

with  $dist = 0.5136$ ,  $\hat{\pi} = 0.4881$ , and the normalized matrices are:

$$
\hat{\beta}_{500}^{*} = \left(\begin{array}{c} 0.0911 & 0.4056 \\ 0.1236 & 0.3226 \\ 0.4158 & 0.0959 \\ 0.3693 & 0.1757 \end{array}\right), \ \ \hat{\gamma}_{500}^{*} = \left(\begin{array}{c} 0.1091 & 0.4010 \\ 0.1474 & 0.2969 \\ 0.2080 & 0.2159 \\ 0.5354 & 0.0859 \end{array}\right),
$$

For the sample size  $T=1000$ , we get the estimated matrices:

$$
\hat{B}_{1000}^{IML} = \hat{B}_{1000}^{AML} * \hat{Q}_{1000} = \begin{pmatrix} 0.0480 & 0.0687 \\ 0.0628 & 0.0750 \\ 0.1379 & 0.0220 \\ 0.1225 & 0.0409 \end{pmatrix} \quad \hat{C}_{1000} = {}^{IML} = \hat{C}_{1000}^{AML} * (\hat{Q}'_{1000})^{-1} = \begin{pmatrix} 0.1140 & 0.3718 \\ 0.0794 & 0.3321 \\ 0.1558 & 0.1774 \\ 0.4670 & 0.1901 \end{pmatrix}, \quad \hat{A}_{1000} = \begin{pmatrix} 0.0310 & 0.0266 & 0.0197 & 0.0355 \\ 0.0351 & 0.0299 & 0.0231 & 0.0436 \\ 0.0240 & 0.0183 & 0.0254 & 0.0687 \\ 0.0292 & 0.0233 & 0.0264 & 0.0650 \end{pmatrix},
$$

with  $dist = 0.4119$ ,  $\hat{\pi} = 0.5777$  and the normalized matrices are:

$$
\hat{\beta}_{1000}^* = \left(\begin{array}{ccc} 0.1292 & 0.3323 \\ 0.1692 & 0.3630 \\ 0.3715 & 0.1067 \\ 0.3299 & 0.1978 \end{array}\right), \ \ \hat{\gamma}_{1000}^* = \left(\begin{array}{ccc} 0.1397 & 0.3470 \\ 0.0973 & 0.3099 \\ 0.1908 & 0.1655 \\ 0.5720 & 0.1774 \end{array}\right),
$$

where  $\widehat{dist}$  is a weighted measure of distance between  $\hat{A}_T$  and  $A_0$  (see online Appendix 3.2).

Additional estimation results are provided in on-line Appendix 3. In on-line Appendix 3.2 we show the unconstrained OLS estimation of matrix A based on the Seemingly Unrelated Regressions (SUR) model. The approach is simple to implement, although much less accurate than the IML method because the constraint of a known rank is not used and the OLS provides negative values of some estimated coefficients, which are incompatible with the parameter set. The results of the AML step of our procedure are given in on-line Appendix 3.3. We observe that the estimators of matrix A converge, while the AML estimators of  $B^*$  and  $C^*$  show much more variation due to the identification issue.

## 7 Concluding Remarks

We introduce the Identifying Maximum Likelihood (IML) method for estimation of the identified set of NMF's and derive the asymptotic distribution of the estimated set and of specific elements of that set. Moreover, we provide a ML estimator of the non-negative matrix  $A_0$  under a given non-negative rank constraint and derive its asymptotic distribution. The proposed approach is related to the nonparametric identification in mixture models. Although the non-negative matrix factorization (NMF) is a well-known technique of dimension reduction for non-negative matrices, it is used wihout an associated probability model for the observed data. Our paper fills this gap by considering a stochastic dynamic network model.

Our approach can be used in a variety of applications with a lack of local identifiability when an element of the identified set is characterized as a solution of auxiliary optimizations, and the identified set can be parametrized given this element considered as a new origin of the identified set (manifold). In practice, the dimension of the parametric identified set can be very large and impossible to be represented in a low-dimensional figure. However, it is possible to illustrate various elements or cuts of that set, which have structural interpretations and are easier to plot and discuss [see on-line Appendix 4].

Partial identifications of the same type are encountered in Singular Value Decomposition (see on-line Appendix 5) and the analysis of monetary shock effects under the sign restrictions on impulse response functions [Baumeister and Hamilton (2015), Ariaz, Rubio-Ramirez and Waggoner (2018), Granzeria et al. (2018)]. For example, in a VAR(1) model  $Y_t = \Phi Y_{t-1} + u_t$ , with  $u_t \sim N(0, \Sigma)$  and  $\Sigma = AA'$ , the matrix A is not identifiable. The identified set can be reduced by imposing a restriction so that a "monetary policy impulse vector (a column of  $A$ ) implies negative responses on prices and non-borrowed reserves and positive responses on federal funds rate, at all horizons". Then, the reduced identified set of these impulse vectors can be easily parametrized [see, a description of the IRF identified set in Gourieroux, Jasiak (2024), Proposition 2. To select the most interpretable  $\beta_k$ 's or  $\gamma'_k s$ , one can use an approach similar to the factor analysis of time series, where the dynamic factors are obtained from the SVD, by projecting factorial directions on the observable time series (the so-called mimicking dynamic factors). A similar approach can be applied in the NMF framework with partial identification.

#### **REFERENCES**

Anderson, T., (1963) : "Asymptotic Theory for Principal Component Analysis", Ann. Math. Statist., 34, 122-148.

Andrews, D. and G. Soarez (2010): "Inference for Parameters Defined by Moment Inequalities Using Generalized Moment Selection", Econometrica, 78, 119-157.

Ariaz, J., Rubio-Ramirez, J. and D. Waggoner (2018): "Inference Based on Structural Vector Autoregressions Identified with Sign and Zero Restrictions: Theory and Applications", Econometrica, 86, 685-720.

Baumeister, C. and J. Hamilton (2015): "Sign Restrictions, Structural Vector Autoregressions and Useful Prior Information", Econometrica, 83, 1963-1999.

Berman, A., and R., Plemmons (1994) : "Nonnegative Matrices in the Mathematical Sciences", Philadelphia, Classics in Applied Mathematics, SIAM, Elsevier.

Blume, L., Brock, W., Durlauf, S., and Y., Ionnides (2011) : "Identification of Social Interactions", Handbook of Social Economics, Vol 1, 853-964.

Bramoulle, Y., Djebbari, H., and B., Fortin (2009) : "Identification of Peer Effects Through Social Networks", Journal of Econometrics, 150, 41-55.

Brie, D. (2015) : "On the Uniqueness and Admissible Solutions of Nonnegative Matrix Factorization", Winter School, Search for Latent Variables : ICA's, Tensors and NMF, Villard de Lans.

Brock, W., and S., Durlauf (2010) : "Adoption Curves and Social Interactions", Journal of the European Economics Association, 8, 235-251.

Cai, J., Yang, D. Zhu, W. Shen, H. and L. Zhao (2021): "Network Regressions and Supervised Centrality Estimation", DP. University of Pennsylvania.

Cameron, C. and P. Trivedi (2008): "Regression Analysis of Count Data", Economeric Society Monographs, 30, Cambridge University Press.

Canay, I., and A. Shaikh (2017): "Practical and Theoretical Advances in Inference for Partially Identified Models", in Advances in Economics and Econometrics, B. Homore, A. Pakes, M. Piazzesi, and L. Samuelson eds., Vol 2., Chap 9, 271-306, Cambridge University Press.

Chen, X. , Christensen, M. and E. Tamer (2018): "Monte-Carlo Confidence Sets for Identified Sets", Econometrica, 86, 1965-2018.

Chen, M., Fernandez-Val, I., and M., Weidner (2021) : "Nonlinear Factor Model for Network and Panel Data", Journal of Econometrics, 220, 296-324.

Chernozhukov, V., Hong, H., and E., Tamer (2007) : "Estimation and Confidence Regions for Parameter Sets in Econometric Models", Econometrica, 75, 1243-1284.

Compiani, G., and Y., Kitamura (2016) : "Using Mixtures in Econometric Models : A Brief Review and Some New Results", Econometrics Journal, 19, 95-127.

Davezies, L., D'Haultfoeuille, X., and L., Laage (2022) : "Identification and Estimation of Average Marginal Effects in Fixed Effects Logit Models", arXiv 2105.00879.

De Paula, A. (2017) : "Econometrics of Network Models", in Advances in Economics and Econometrics, Eds, Honore, B., Pakes, A., Piazzesi, M, and L., Samuelson, Cambridge Univ., Press Chapter 8.

Del Buono, N., Esposito, F. and N. Gillis (2019) "Orthogonal Joint Sparse NMF for Microarray Data Analysis", Journal of Math. Biology, 79, 223.

Donnet, S. and S. Robin (2021): "Accelerating Bayesian Estimation for Network Poisson Models Using Frequentist Variational Estimates", JRSS Series C, 70, 858- 885.

Engle, R. and J. Rangel (2008): "The Spline GARCH Model for Unconditional Volatility and its Global Macroeconomic Causes", Review of Financial Studies, 21, 1187-1222.

Esposito, F. (2021): "A Review of Initialization Methods for Nonnegative Matrix Factorization: Towards Omics Data Experiments", Mathematics, 9, 10060.

Fahrenwaldt, M., Weber, S., and K., Weske (2018) : "Pricing of Cyber Insurance Contracts in a Network Model", ASTIN Bulletin, 48, 1175-1218.

Fernandez-Val, I., and M., Weidner (2016) : "Individual and Time Effects in Nonlinear Panel Models with Large  $N, T$ ", Journal of Econometrics 192, 291-312.

Fokianos, K. (2022): "Multivariate Count Time Series Modelling", Econometrics and Statistics, forthcoming.

Gill, P., Murray, W., and M., Saunders (2002) : "SNOPT : An SQP Algorithm for Large Scale Constrained Optimization", SIAM J. Optim., K, 979-1006.

Gillis, N. (2020): "Nonnegative Matrix Factorization", SIAM , Philadelphia.

Gourieroux, C., and J., Jasiak (2024) : "Statistical Inference for SVAR Model with Sign Restriction", CREST DP.

Gourieroux, C. and Y., Lu (2019): "Negative Binomial Autoregressive Process with Stochastic Intensity", Journal of Time Series Analysis, 40, 225-247.

Gourieroux, C. and Y., Lu (2023): "Susceptible-Infected- Recovered Model with Stochastic Transmission", Canadian Journal of Statistics, forthcoming.

Gourieroux, C., and A., Monfort (1995) : "Statistics and Econometric Models", Vol 1, Cambridge University Press.

Gourieroux, C., Monfort, A., and E., Renault (1990) : "Biilinear Constraints: Estimation and Test", Essays in Honor of Edmond Malinvaud, Empirical Economics, MIT Press, 166-191.

Granzeria, E., Moon, H. and F. Schorfheide (2018): "Inference for VARs Identified with Sign Restrictions", Quantitative Economics, 9, 1987-2111.

Grippo, L. and M. Sciandrone (2002): "On the Convergence of Block Nonlinear Gauss Seidel Method under Convex Constraints", Optimization Research Letters, 26, 127- 136.

Hafner, C. and O. Linton (2010): "Efficient Estimation of a Multivariate Multiplicative Volatility Model", Journal of Econometrics, 159, 55-73.

Hall, P., and X., Zhou (2003) : "Nonparametric Estimation of Component Distributions in a Multivariate Mixture", Annals of Statistics, 31, 201-224.

Hastie, T., Mazunder, M., Lee, J. and R., Zadeh (2014): "Matrix Completion and Low Rank SVD via Fast Alternating Least Squares", Journal of Machine Learning Research, 16, 3367-3402.

Hautsch, N., Ohrkin, O., and R. Alexamder (2023): "Maximum Likelihood Estimation Under the Zig-Zag Algorithm", Journal of Financial Econometrics, 21, 1346-1375. Henry, M., Kitamura, Y., and B., Salanie (2014) : "Partial Identification of Finite Mixtures in Econometric Models", Quantitative Economics, 5, 123-144.

Hien, L., and N. Gillis (2021): "Algorithm for Nonnegative Matrix Factorization with the Kullblack-Leibler Divergence", Arxiv 2010-01935V2.

Imbens, G., and C., Manski (2004) : "Confidence Intervals for Partially Identified Parameters", Econometrica, 72, 1845-1857.

Kasahara, H., and K., Shimotsu (2009) : "Nonparametric Identification of Finite Mixture Models of Dynamic Discrete Choices", Econometrica, 77, 135-175.

Kim, H., and H., Park (2008) : "Nonnegative Matrix Factorization Based on Alternating Nonnegativity Constrained Least Squares and Active Set Method", SIAM Journal of Matrix Anal. Appl., 30, 716-730.

Kline, B. and E. Tamer (2023): "Recent Developments in Partial Identification", Annual Review of Economics, 15, 125-150.

Laurberg, H., Christensen, M., Plumbey, M., Hansen, L., and S., Jensen (2008) : "Theorems on Positive Data : On the Uniqueness of NMF", Computational Intelligence and Neuroscience, ID 764206.

Lee, L., Liu, X., and S., Lin (2010) : "Specification and Estimation of Social Interaction Models with Network Structure, Contextual Factors, Correlation and Fixed Effects", Econometrics Journal, 13, 145-176.

Lee, D., and H., Seung (1999) : "Learning the Parts of Objects by Nonnegative Matrix Factorization", Nature, 401, 788-791.

Liu, X. (2005) : "Global Convergence on an Active Set SQP for Inequality Constrained Optimization", Journal of Computational and Applied Mathematics, 180, 201-211.

Liu, L., Poirier, A., and J., Shin (2023) :" Identification and Estimation of Average Partial Effects in Semi-Parametric Binary Response Panel Models", arXiv : 2105.12891.

Lu, Y., Zhang, J. and W. Zhu (2024): "Cyber Risk Modeling: a Discrete Multivariate Count Process Approach", Scandinavian Actuarial Journal, 1-31.

Manski, C. (1993) : "Identification of Endogenous Social Effects : The Reflection Problem", The Review of Economic Studies, 60, 531-542.

Molchanov, I. and F. Molinari (2018): "Random Sets in Econometrics", Cambridge University Press.

Paatero, P., and U., Tapper (1994): "Positive Matrix Factorization: A Nonnegative Factor Model with Optimal Utilisation of Error Estimates of Data Values", Environmetrics, 5, 111-126.

Schachtner, R., Pappel, G., and E., Lang (2011) : "Toward Unique Solutions of Nonnegative Matrix Factorization Problems by a Determinant Criterion", Digital Signal Processing, 21, 528-534.

Shi, X. and M. Shum (2015): "Simple Two-Stage Inference for a Class of Partially Identified Models", Econometric Theory, 31, 493-520.

Stoye, J. (2009) : "More on Confidence Intervals for Partially Identified Parameters", Econometrica, 77, 1299-1315.

Tipping, H., and C., Bishop (1999) : "Probabilistic Principal Component Analysis", JRSS, B, 61, 611-622.

Vrahalis, M., Magoulas, G. and V. Piagianakos (2003): "From Linear to Nonlinear Iterative Methods", Applied Numerical Mathematics, 45, 59-77.

## Appendix 1: Proofs of Propositions 5 and 6

#### 1. Proof or Proposition 5

i) The proof is standard and based on the asymptotic expansion of the first-order conditions on the Lagrangean for the equality restrictions (note that the inequality restrictions are not binding if  $\alpha_0^*$  belongs in the interior of  $\mathcal{A}_0$ ). These asymptotic expansions are :

$$
\begin{pmatrix} J_0 & D \ D' & 0 \end{pmatrix} \begin{bmatrix} \sqrt{T}(\hat{\alpha}_T - \alpha_0^*) \\ \sqrt{T}(\hat{\lambda}_T - \lambda_0^*) \end{bmatrix} \simeq \begin{bmatrix} \frac{1}{\sqrt{T}} \sum_{t=1}^T \frac{\partial \log l}{\partial \alpha}(y_t | y_{t-1}; \alpha_0^*) \\ 0 \end{bmatrix}, \quad (a.1)
$$

where  $\hat{\lambda}_T$  is the associated estimator of the Lagrange multipliers. Because  $\alpha_0^*$  is a maximizer of  $E_0 \log l(y_t | y_{t-1}, \alpha)$ , the score in (4.13) is asymptotically normally distributed with mean zero and variance J.

Then, the result follows whenever the matrix  $\begin{pmatrix} J_0 & D \\ D_0 & 0 \end{pmatrix}$  $D'$  0 ) is invertible.  $23$ 

ii) Let us now discuss this invertibility condition by finding the null space of this matrix, i.e. the solutions  $\theta$ ,  $\lambda$  of the system :  $J_0\theta + D\lambda = 0$ ,  $D'\theta = 0$ .

We know that  $D'J_0\theta + D'D\lambda = 0 \Rightarrow \lambda = -(D'D)^{-1}D'J_0\theta$ . Then the system in  $\theta$  only is :  $(Id - P)J_0\theta = 0$ ,  $D'\theta = 0$ , where P is the orthogonal projector on the space generated by D.

Since the columns of  $(Id - P)J_0$  are orthogonal to the columns of D, we see that  $\theta = 0$  is the unique solution of the above system, if and only if :  $Rk((Id - P)J_0) =$  $\dim \alpha - \dim q - 1 - 2K$ . This statement is Assumption a.2 viii) in on-line Appendix 2.

#### 2. Proof or Proposition 6

The asymptotic first-order conditions involve  $J_0$ √  $\overline{T}(\hat{\theta}_T - \alpha_0^*) + D_1$  $\frac{1}{T}$  asymptotic first-order conditions involve  $J_0\sqrt{T}(\hat{\theta}_T - \alpha_0^*) + D_1\sqrt{T}(\hat{\lambda}_T - \lambda_{10}^*) +$  $D_2\sqrt{T}(\hat{\lambda}_{2T}-\lambda_{20}^*)$  and  $D_1'\sqrt{T}(\hat{\alpha}_T-\alpha_0^*)+D_2'\sqrt{T}(\hat{\alpha}_T-\alpha_0^*)$  in the left hand side of system (4.13). Asymptotically, a change of the benchmark modifies the matrix  $D_2$  as well as the associated Lagrange multipliers by linear transformations R and  $R^{-1}$ , respectively. Then, the first-order conditions provide the same solution for  $\sqrt{T}(\hat{\alpha}_T - \alpha_0^*)$  when  $D_2$ is replaced by  $\tilde{D}_2 = D_0 R$  and  $\hat{\lambda}_{2T} - \lambda_{20}^*$  by  $\hat{\tilde{\lambda}}_{2T} - \tilde{\lambda}_{20}^* = R^{-1}(\hat{\lambda}_{2T} - \lambda_0^*)$ , where R is invertible. This proves that the asymptotic variance-covariance matrix is independent of the choice of the additional optimization criterion.

<sup>&</sup>lt;sup>23</sup>We cannot use the usual block formula to compute  $J_0^{11}$  [see e.g. Gourieroux and Monfort, Section 10.3.b)], because  $J_0$  is not invertible due to the identification issue. However, it is easy to check that a closed form expression of the asymptotic variance of the estimator is :  $[(Id-P)J_0(Id-P)J_1(Id-P)J_2(Id-P)J_2(Id-P)J_3(Id-P)J_3(Id-P)J_4(Id-P)J_5(Id-P)J_6(Id-P)J_7(Id-P)J_8(Id-P)J_8(Id-P)J_9(Id-P)J_9(Id-P)J_9(Id-P)J_8(Id-P)J_9(Id-P)J_9(Id-P)J_9(Id$  $(P) + P^{-1}(Id - P)J_0(Id - P)[(Id - P)J_0(Id - P) + P]^{-1}$ , where P is the orthogonal projector on the space generated by D.