



Invited article

Independent component analysis: algorithms and applications

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Abstract

A fundamental problem in neural network research, as well as in many other disciplines, is finding a suitable representation of multivariate data, i.e. random vectors. For reasons of computational and conceptual simplicity, the representation is often sought as a linear transformation of the original data. In other words, each component of the representation is a linear combination of the original variables. Well-known linear transformation methods include principal component analysis, factor analysis, and projection pursuit. Independent component analysis (ICA) is a recently developed method in which the goal is to find a linear representation of non-Gaussian data so that the components are statistically independent, or as independent as possible. Such a representation seems to capture the essential structure of the data in many applications, including feature extraction and signal separation. In this paper, we present the basic theory and applications of ICA, and our recent work on the subject. © 2000 Published by Elsevier Science Ltd.

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1. Motivation

Imagine that you are in a room where two people are speaking simultaneously. You have two microphones, which you hold in different locations. The microphones give you two recorded time signals, which we could denote by $x_1(t)$ and $x_2(t)$, with x_1 and x_2 the amplitudes, and t the time index. Each of these recorded signals is a weighted sum of the speech signals emitted by the two speakers, which we denote by $s_1(t)$ and $s_2(t)$. We could express this as a linear equation:

$$x_1(t) = a_{11}s_1 + a_{12}s_2 \quad (1)$$

$$x_2(t) = a_{21}s_1 + a_{22}s_2 \quad (2)$$

where a_{11} , a_{12} , a_{21} , and a_{22} are some parameters that depend on the distances of the microphones from the speakers. It would be very useful if you could now estimate the two original speech signals $s_1(t)$ and $s_2(t)$, using only the recorded signals $x_1(t)$ and $x_2(t)$. This is called the *cocktail-party problem*. For the time being, we omit any time delays or other extra factors from our simplified mixing model.

As an illustration, consider the waveforms in Figs. 1 and 2. These are, of course, not realistic speech signals, but suffice for this illustration. The original speech signals

could look something like those in Fig. 1 and the mixed signals could look like those in Fig. 2. The problem is to recover the data in Fig. 1 using only the data in Fig. 2.

Actually, if we knew the parameters a_{ij} , we could solve the linear equation in (1) by classical methods. The point is, however, that if you do not know the a_{ij} , the problem is considerably more difficult.

One approach to solving this problem would be to use some information on the statistical properties of the signals $s_i(t)$ to estimate the a_{ij} . Actually, and perhaps surprisingly, it turns out that it is enough to assume that $s_1(t)$ and $s_2(t)$, at each time instant t , are *statistically independent*. This is not an unrealistic assumption in many cases, and it need not be exactly true in practice. The recently developed technique of Independent Component Analysis, or ICA, can be used to estimate the a_{ij} based on the information of their independence, which allows us to separate the two original source signals $s_1(t)$ and $s_2(t)$ from their mixtures $x_1(t)$ and $x_2(t)$. Fig. 3 gives the two signals estimated by the ICA method. As can be seen, these are very close to the original source signals (their signs are reversed, but this has no significance).

Independent component analysis was originally developed to deal with problems that are closely related to the cocktail-party problem. Since the recent increase of interest in ICA, it has become clear that this principle has a lot of other interesting applications as well.

Consider, for example, electrical recordings of brain activity as given by an electroencephalogram (EEG). The EEG data consists of recordings of electrical potentials in

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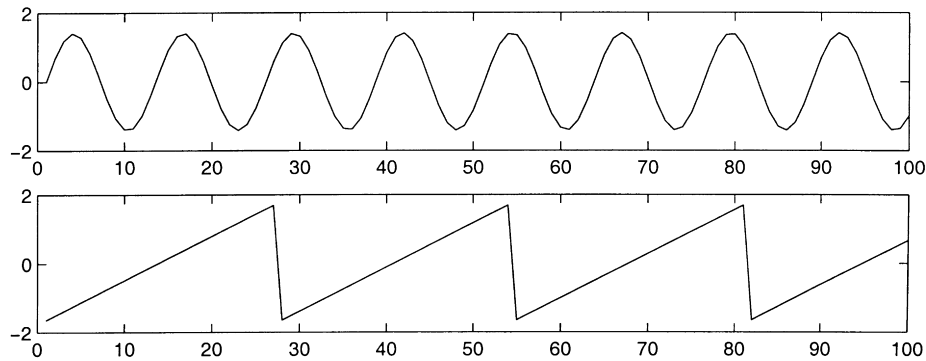


Fig. 1. The original signals.

many different locations on the scalp. These potentials are presumably generated by mixing some underlying components of brain activity. This situation is quite similar to the cocktail-party problem: we would like to find the original components of brain activity, but we can only observe mixtures of the components. ICA can reveal interesting information on brain activity by giving access to its independent components.

Another, very different application of ICA is on feature extraction. A fundamental problem in digital signal processing is to find suitable representations for image, audio or other kind of data for tasks like compression and denoising. Data representations are often based on (discrete) linear transformations. Standard linear transformations widely used in image processing are the Fourier, Haar, cosine transforms etc. Each of them has its own favorable properties (Gonzalez & Wintz, 1987).

It would be most useful to estimate the linear transformation from the data itself, in which case the transform could be ideally adapted to the kind of data that is being processed. Fig. 4 shows the basis functions obtained by ICA from patches of natural images. Each image window in the set of training images would be a superposition of these windows so that the coefficient in the superposition is independent. Feature extraction by ICA will be explained in more detail later on.

All of the applications described above can actually be formulated in a unified mathematical framework, that of

ICA. This is a very general-purpose method of signal processing and data analysis.

In this review, we cover the definition and underlying principles of ICA in Sections 2 and 3. Then, starting from Section 4, the ICA problem is solved on the basis of minimizing or maximizing certain contrast functions; this transforms the ICA problem to a numerical optimization problem. Many contrast functions are given and the relations between them are clarified. Section 5 covers a useful preprocessing that greatly helps solving the ICA problem, and Section 6 reviews one of the most efficient practical learning rules for solving the problem, the FastICA algorithm. Then, in Section 7, typical applications of ICA are covered: removing artifacts from brain signal recordings, finding hidden factors in financial time series, and reducing noise in natural images. Section 8 concludes the text.

2. Independent component analysis

2.1. Definition of ICA

To rigorously define ICA (Comon, 1994; Jutten & Herault, 1991), we can use a statistical “latent variables” model. Assume that we observe n linear mixtures x_1, \dots, x_n of n independent components

$$x_j = a_{j1}s_1 + a_{j2}s_2 + \dots + a_{jn}s_n, \quad \text{for all } j. \quad (3)$$

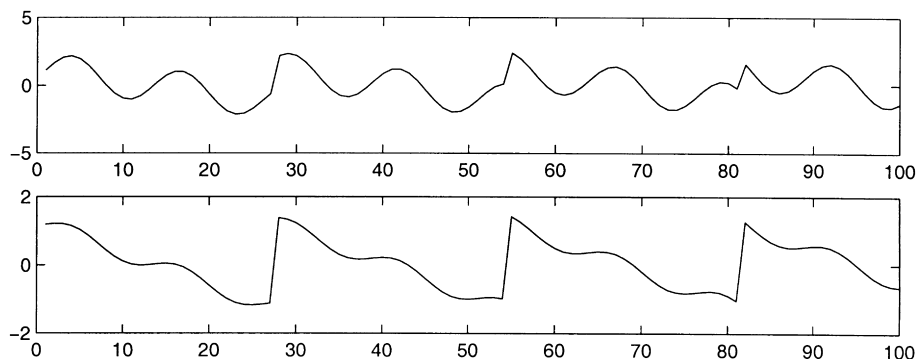


Fig. 2. The observed mixtures of the source signals in Fig. 1.

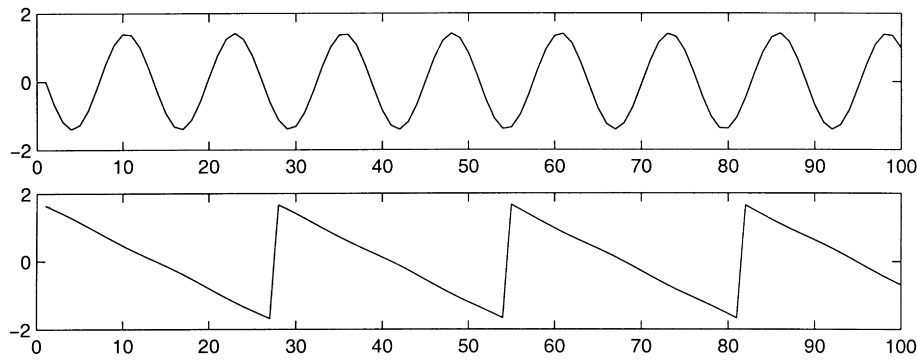


Fig. 3. The estimates of the original source signals, estimated using only the observed signals in Fig. 2. The original signals were very accurately estimated, up to multiplicative signs.

We have now dropped the time index t ; in the ICA model, we assume that each mixture x_j as well as each independent component s_k is a random variable, instead of a proper time signal. The observed values $x_j(t)$, e.g. the microphone signals in the cocktail party problem, are then a sample of this random variable. Without loss of generality, we can assume that both the mixture variables and the independent components have zero mean: If this is not true, then the

observable variables x_i can always be centered by subtracting the sample mean, which makes the model zero-mean.

It is convenient to use vector–matrix notation instead of the sums like in the previous equation. Let us denote by \mathbf{x} the random vector whose elements are the mixtures x_1, \dots, x_n , and likewise by \mathbf{s} the random vector with elements s_1, \dots, s_n . Let us denote by \mathbf{A} the matrix with elements a_{ij} . Generally, bold lower case letters indicate vectors and bold upper-case

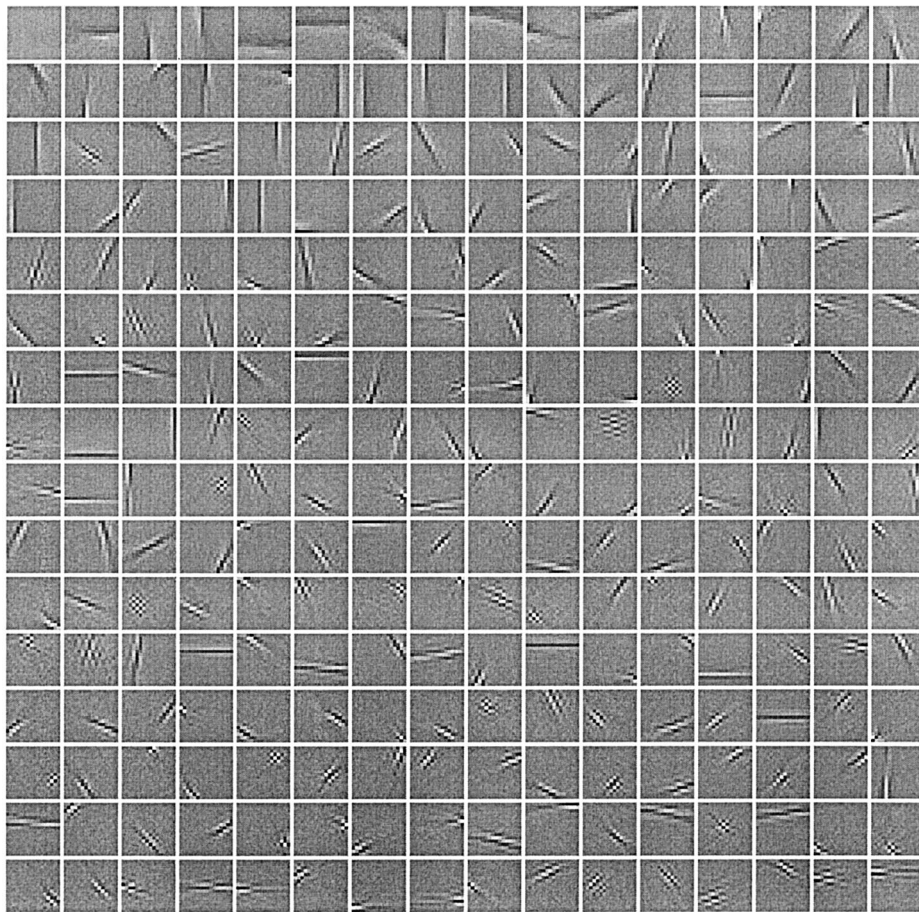


Fig. 4. Basis functions in ICA of natural images. The input window size was 16×16 pixels. These basis functions can be considered as the independent features of images.

letters denote matrices. All vectors are understood as column vectors; thus \mathbf{x}^T , or the transpose of \mathbf{x} , is a row vector. Using this vector–matrix notation, the above mixing model is written as

$$\mathbf{x} = \mathbf{A}\mathbf{s}. \quad (4)$$

Sometimes we need the columns of matrix \mathbf{A} ; denoting them by \mathbf{a}_j the model can also be written as

$$\mathbf{x} = \sum_{i=1}^n \mathbf{a}_i s_i. \quad (5)$$

The statistical model in Eq. (4) is called independent component analysis, or ICA model. The ICA model is a generative model, which means that it describes how the observed data are generated by a process of mixing the components s_i . The independent components are latent variables, meaning that they cannot be directly observed. Also the mixing matrix is assumed to be unknown. All we observe is the random vector \mathbf{x} , and we must estimate both \mathbf{A} and \mathbf{s} using it. This must be done under as general assumptions as possible.

The starting point for ICA is the very simple assumption that the components s_i are statistically *independent*. Statistical independence will be rigorously defined in Section 3. It will be seen below that we must also assume that the independent component must have *non-Gaussian* distributions. However, in the basic model we do *not* assume these distributions known (if they are known, the problem is considerably simplified). For simplicity, we are also assuming that the unknown mixing matrix is square, but this assumption can be sometimes relaxed, as explained in Section 4.5. Then, after estimating the matrix \mathbf{A} , we can compute its inverse, say \mathbf{W} , and obtain the independent component simply by:

$$\mathbf{s} = \mathbf{W}\mathbf{x}. \quad (6)$$

ICA is very closely related to the method called *blind source separation* (BSS) or blind signal separation. A “source” means here an original signal, i.e. independent component, like the speaker in a cocktail party problem. “Blind” means that we know very little, if anything, on the mixing matrix, and make little assumptions on the source signals. ICA is one method, perhaps the most widely used, for performing blind source separation.

In many applications, it would be more realistic to assume that there is some noise in the measurements (Hyvärinen, 1998a; Hyvärinen, 1999c), which would mean adding a noise term in the model. For simplicity, we omit any noise terms, since the estimation of the noise-free model is difficult enough in itself, and seems to be sufficient for many applications.

2.2. Ambiguities of ICA

In the ICA model in Eq. (4), it is easy to see that the

following ambiguities will hold:

1. We cannot determine the variances (energies) of the independent components.

The reason is that, both \mathbf{s} and \mathbf{A} being unknown, any scalar multiplier in one of the sources s_i could always be cancelled by dividing the corresponding column \mathbf{a}_i of \mathbf{A} by the same scalar; see Eq. (5). As a consequence, we may quite as well fix the magnitudes of the independent components; as they are random variables, the most natural way to do this is to assume that each has unit variance: $E\{s_i^2\} = 1$. Then the matrix \mathbf{A} will be adapted in the ICA solution methods to take into account this restriction. Note that this still leaves the ambiguity of the sign: we could multiply the independent component by -1 without affecting the model. This ambiguity is, fortunately, insignificant in most applications.

2. We cannot determine the order of the independent components.

The reason is that, again both \mathbf{s} and \mathbf{A} being unknown, we can freely change the order of the terms in the sum in Eq. (5), and call any of the independent components the first one. Formally, a permutation matrix \mathbf{P} and its inverse can be substituted in the model to give $\mathbf{x} = \mathbf{A}\mathbf{P}^{-1}\mathbf{P}\mathbf{s}$. The elements of $\mathbf{P}\mathbf{s}$ are the original independent variables s_j , but in another order. The matrix $\mathbf{A}\mathbf{P}^{-1}$ is just a new unknown mixing matrix, to be solved by the ICA algorithms.

2.3. Illustration of ICA

To illustrate the ICA model in statistical terms, consider two independent components that have the following uniform distributions:

$$p(s_i) = \begin{cases} \frac{1}{2\sqrt{3}} & \text{if } |s_i| \leq \sqrt{3} \\ 0 & \text{otherwise} \end{cases} \quad (7)$$

The range of values for this uniform distribution were chosen so as to make the mean zero and the variance equal to one, as was agreed in the previous section. The joint density of s_1 and s_2 is then uniform on a square. This follows from the basic definition that the joint density of two independent variables is just the product of their marginal densities (see Eq. (10)): we need to simply compute the product. The joint density is illustrated in Fig. 5 by showing data points randomly drawn from this distribution.

Now let us mix these two independent components. Let us take the following mixing matrix:

$$\mathbf{A}_0 = \begin{pmatrix} 2 & 3 \\ 2 & 1 \end{pmatrix} \quad (8)$$

This gives us two mixed variables, x_1 and x_2 . It is easily computed that the mixed data has a uniform distribution on a parallelogram, as shown in Fig. 6. Note that the random

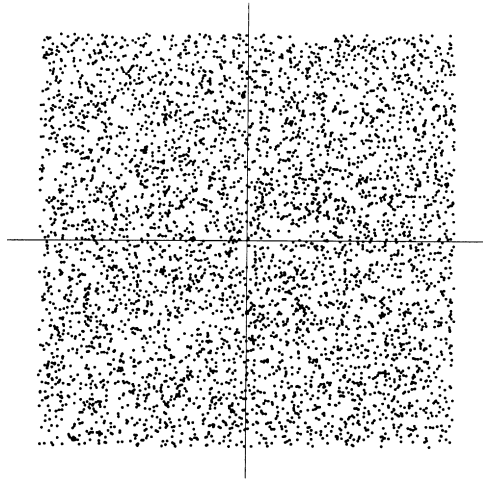


Fig. 5. The joint distribution of the independent components s_1 and s_2 with uniform distributions. Horizontal axis: s_1 , vertical axis: s_2 .

variables x_1 and x_2 are not independent any more; an easy way to see this is to consider, whether it is possible to predict the value of one of them, say x_2 , from the value of the other. Clearly if x_1 attains one of its maximum or minimum values, then this completely determines the value of x_2 . They are therefore not independent. (For variables s_1 and s_2 the situation is different: from Fig. 5 it can be seen that knowing the value of s_1 does not in any way help in guessing the value of s_2 .)

The problem of estimating the data model of ICA is now to estimate the mixing matrix \mathbf{A}_0 using only information contained in the mixtures x_1 and x_2 . Actually, you can see from Fig. 6 an intuitive way of estimating \mathbf{A} : the edges of the parallelogram are in the directions of the columns of \mathbf{A} . This means that we could, in principle, estimate the ICA model by first estimating the joint density of x_1 and x_2 , and then locating the edges. So, the problem seems to have a solution.

In reality, however, this would be a very poor method because it only works with variables that have exactly uniform distributions. Moreover, it would be computationally quite complicated. What we need is a method that works for any distributions of the independent components, and works fast and reliably.

Next we shall consider the exact definition of independence before starting to develop methods for estimation of the ICA model.

3. What is independence?

3.1. Definition and fundamental properties

To define the concept of independence, consider two scalar-valued random variables y_1 and y_2 . Basically, the variables y_1 and y_2 are said to be independent if information on the value of y_1 does not give any information on the value

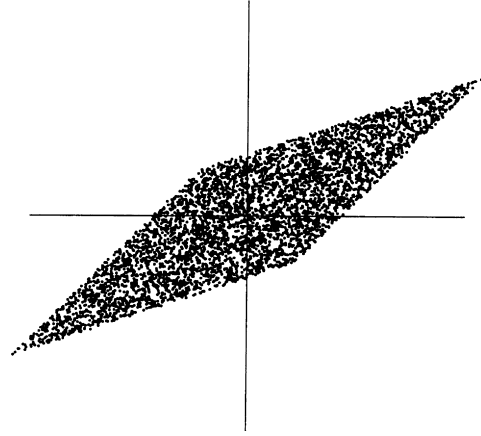


Fig. 6. The joint distribution of the observed mixtures x_1 and x_2 . Horizontal axis: x_1 , vertical axis: x_2 .

of y_2 , and vice versa. Above, we noted that this is the case with the variables s_1, s_2 but not with the mixture variables x_1, x_2 .

Technically, independence can be defined by the probability densities. Let us denote by $p(y_1, y_2)$ the joint probability density function (pdf) of y_1 and y_2 . Let us further denote by $p_1(y_1)$ the marginal pdf of y_1 , i.e. the pdf of y_1 when it is considered alone:

$$p_1(y_1) = \int p(y_1, y_2) dy_2, \quad (9)$$

and similarly for y_2 . Then we define that y_1 and y_2 are independent if and only if the joint pdf is *factorizable* in the following way:

$$p(y_1, y_2) = p_1(y_1)p_2(y_2). \quad (10)$$

This definition extends naturally for any number n of random variables, in which case the joint density must be a product of n terms.

The definition can be used to derive a most important property of independent random variables. Given two functions, h_1 and h_2 , we always have

$$E\{h_1(y_1)h_2(y_2)\} = E\{h_1(y_1)\}E\{h_2(y_2)\}. \quad (11)$$

This can be proven as follows:

$$\begin{aligned} E\{h_1(y_1)h_2(y_2)\} &= \iint h_1(y_1)h_2(y_2)p(y_1, y_2)dy_1 dy_2 \\ &= \iint h_1(y_1)p_1(y_1)h_2(y_2)p_2(y_2)dy_1 dy_2 \\ &= \int h_1(y_1)p_1(y_1)dy_1 \int h_2(y_2)p_2(y_2)dy_2 \\ &= E\{h_1(y_1)\}E\{h_2(y_2)\}. \end{aligned} \quad (12)$$

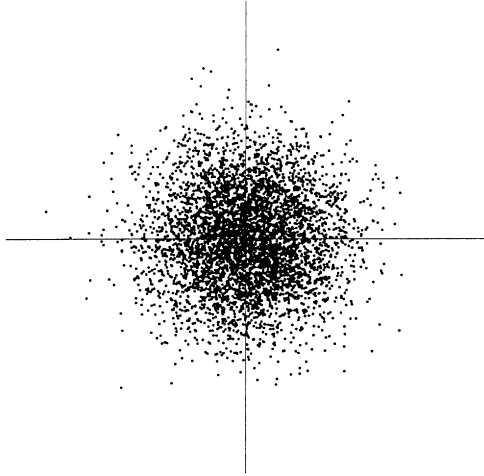


Fig. 7. The multivariate distribution of two independent Gaussian variables.

3.2. Uncorrelated variables are only partly independent

A weaker form of independence is uncorrelatedness. Two random variables y_1 and y_2 are said to be uncorrelated, if their covariance is zero:

$$E\{y_1 y_2\} - E\{y_1\}E\{y_2\} = 0 \quad (13)$$

If the variables are independent, they are uncorrelated, which follows directly from Eq. (11), taking $h_1(y_1) = y_1$ and $h_2(y_2) = y_2$.

On the other hand, uncorrelatedness does *not* imply independence. For example, assume that (y_1, y_2) are discrete valued and follow such a distribution that the pair are with probability 1/4 equal to any of the following values: $(0,1), (0,-1), (1,0), (-1,0)$. Then y_1 and y_2 are uncorrelated, as can be simply calculated. On the other hand,

$$E\{y_1^2 y_2^2\} = 0 \neq \frac{1}{4} = E\{y_1^2\}E\{y_2^2\}. \quad (14)$$

so the condition in Eq. (11) is violated, and the variables cannot be independent.

Since independence implies uncorrelatedness, many ICA methods constrain the estimation procedure so that it always gives uncorrelated estimates of the independent components. This reduces the number of free parameters, and simplifies the problem.

3.3. Why Gaussian variables are forbidden

The fundamental restriction in ICA is that the independent components must be non-Gaussian for ICA to be possible.

To see why Gaussian variables make ICA impossible, assume that the mixing matrix is orthogonal and the s_i are Gaussian. Then x_1 and x_2 are Gaussian, uncorrelated, and of unit variance. Their joint density is given by

$$p(x_1, x_2) = \frac{1}{2\pi} \exp\left(-\frac{x_1^2 + x_2^2}{2}\right) \quad (15)$$

This distribution is illustrated in Fig. 7. The figure shows that the density is completely symmetric. Therefore, it does not contain any information on the directions of the columns of the mixing matrix \mathbf{A} . This is why \mathbf{A} cannot be estimated.

More rigorously, one can prove that the distribution of any orthogonal transformation of the Gaussian (x_1, x_2) has exactly the same distribution as (x_1, x_2) , and that x_1 and x_2 are independent. Thus, in the case of Gaussian variables, we can only estimate the ICA model up to an orthogonal transformation. In other words, the matrix \mathbf{A} is not identifiable for Gaussian independent components. (Actually, if just one of the independent components is Gaussian, the ICA model can still be estimated.)

4. Principles of ICA estimation

4.1. “Non-Gaussian is independent”

Intuitively speaking, the key to estimating the ICA model is non-Gaussianity. Actually, without non-Gaussianity the estimation is not possible at all, as mentioned in Section 3.3. This is at the same time probably the main reason for the rather late resurgence of ICA research: In most of classical statistical theory, random variables are assumed to have Gaussian distributions, thus precluding any methods related to ICA.

The Central Limit Theorem, a classical result in probability theory, tells that the distribution of a sum of independent random variables tends toward a Gaussian distribution, under certain conditions. Thus, a sum of two independent random variables usually has a distribution that is closer to Gaussian than any of the two original random variables.

Let us now assume that the data vector \mathbf{x} is distributed according to the ICA data model in Eq. (4), i.e. it is a mixture of independent components. For simplicity, let us assume in this section that all the independent components have identical distributions. To estimate one of the independent components, we consider a linear combination of the x_i (see Eq. (6)); let us denote this by $y = \mathbf{w}^T \mathbf{x} = \sum_i w_i x_i$, where \mathbf{w} is a vector to be determined. If \mathbf{w} were one of the rows of the inverse of \mathbf{A} , this linear combination would actually equal one of the independent components. The question is now: How could we use the Central Limit Theorem to determine \mathbf{w} so that it would equal one of the rows of the inverse of \mathbf{A} ? In practice, we cannot determine such a \mathbf{w} exactly, because we have no knowledge of matrix \mathbf{A} , but we can find an estimator that gives a good approximation.

To see how this leads to the basic principle of ICA estimation, let us make a change of variables, defining $\mathbf{z} = \mathbf{A}^T \mathbf{w}$. Then we have $y = \mathbf{w}^T \mathbf{x} = \mathbf{w}^T \mathbf{A} \mathbf{s} = \mathbf{z}^T \mathbf{s}$. y is thus a linear combination of s_i , with weights given by z_i . Since a sum of even two independent random variables is more Gaussian than the original variables, $\mathbf{z}^T \mathbf{s}$ is more Gaussian than any of the s_i and becomes least Gaussian when it in fact

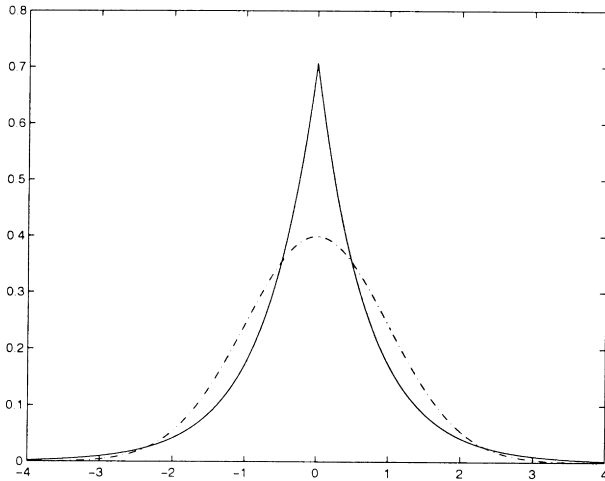


Fig. 8. The density function of the Laplace distribution, which is a typical super-Gaussian distribution. For comparison, the Gaussian density is given by a dashed line. Both densities are normalized to unit variance.

equals one of the s_i . In this case, obviously only one of the elements z_i of \mathbf{z} is non-zero. (Note that the s_i were here assumed to have identical distributions.)

Therefore, we could take as \mathbf{w} a vector that maximizes the non-Gaussianity of $\mathbf{w}^T \mathbf{x}$. Such a vector would necessarily correspond (in the transformed coordinate system) to a \mathbf{z} which has only one non-zero component. This means that $\mathbf{w}^T \mathbf{x} = \mathbf{z}^T \mathbf{s}$ equals one of the independent components!

Maximizing the non-Gaussianity of $\mathbf{w}^T \mathbf{x}$ thus gives us one of the independent components. In fact, the optimization landscape for non-Gaussianity in the n -dimensional space of vectors \mathbf{w} has $2n$ local maxima, two for each independent component, corresponding to s_i and $-s_i$ (recall that the independent components can be estimated only up to a multiplicative sign). To find several independent components, we need to find all these local maxima. This is not difficult, because the different independent components are uncorrelated: We can always constrain the search to the space that gives estimates uncorrelated with the previous ones. This corresponds to orthogonalization in a suitably transformed (i.e. whitened) space.

Our approach here is rather heuristic, but it will be seen in Sections 4.2 and 4.3 that it has a perfectly rigorous justification.

4.2. Measures of non-Gaussianity

To use non-Gaussianity in ICA estimation, we must have a quantitative measure of non-Gaussianity of a random variable, say y . To simplify things, let us assume that y is centered (zero-mean) and has variance equal to one. Actually, one of the functions of preprocessing in ICA algorithms, to be covered in Section 5, is to make this simplification possible.

4.2.1. Kurtosis

The classical measure of non-Gaussianity is kurtosis or the fourth-order cumulant. The kurtosis of y is classically defined by

$$\text{kurt}(y) = E\{y^4\} - 3(E\{y^2\})^2 \quad (16)$$

Actually, since we assumed that y is of unit variance, the right-hand side simplifies to $E\{y^4\} - 3$. This shows that kurtosis is simply a normalized version of the fourth moment $E\{y^4\}$. For a Gaussian y , the fourth moment equals $3(E\{y^2\})^2$. Thus, kurtosis is zero for a Gaussian random variable. For most (but not quite all) non-Gaussian random variables, kurtosis is non-zero.

Kurtosis can be either positive or negative. Random variables that have a negative kurtosis are called sub-Gaussian, and those with positive kurtosis are called super-Gaussian. In statistical literature, the corresponding expressions, platykurtic and leptokurtic, are also used. Super-Gaussian random variables have typically a “spiky” pdf with heavy tails, i.e. the pdf is relatively large at zero and at large values of the variable, while being small for intermediate values. A typical example is the Laplace distribution, whose pdf (normalized to unit variance) is given by

$$p(y) = \frac{1}{\sqrt{2}} \exp(\sqrt{2}|y|) \quad (17)$$

This pdf is illustrated in Fig. 8. Sub-Gaussian random variables, on the other hand, have typically a “flat” pdf, which is rather constant near zero, and very small for larger values of the variable. A typical example is the uniform distribution in Eq. (7).

Typically non-Gaussianity is measured by the absolute value of kurtosis. The square of kurtosis can also be used. These are zero for a Gaussian variable, and greater than zero for most non-Gaussian random variables. There are non-Gaussian random variables that have zero kurtosis, but they can be considered as very rare.

Kurtosis, or rather its absolute value, has been widely used as a measure of non-Gaussianity in ICA and related fields. The main reason is its simplicity, both computational and theoretical. Computationally, kurtosis can be estimated simply by using the fourth moment of the sample data. Theoretical analysis is simplified because of the following linearity property: If x_1 and x_2 are two independent random variables, it holds

$$\text{kurt}(x_1 + x_2) = \text{kurt}(x_1) + \text{kurt}(x_2) \quad (18)$$

and

$$\text{kurt}(\alpha x_1) = \alpha^4 \text{kurt}(x_1) \quad (19)$$

where α is a scalar. These properties can be easily proven using the definition.

To illustrate in a simple example what the optimization landscape for kurtosis looks like, and how independent components could be found by kurtosis minimization or maximization, let us look at a 2-dimensional model $\mathbf{x} = \mathbf{A}\mathbf{s}$.

Assume that the independent components s_1, s_2 have kurtosis values $\text{kurt}(s_1), \text{kurt}(s_2)$, respectively, both different from zero. Remember that we assumed that they have unit variances. We seek for one of the independent components as $y = \mathbf{w}^T \mathbf{x}$.

Let us again make the transformation $\mathbf{z} = \mathbf{A}^T \mathbf{w}$. Then we have $y = \mathbf{w}^T \mathbf{x} = \mathbf{w}^T \mathbf{A} \mathbf{s} = \mathbf{z}^T \mathbf{s} = z_1 s_1 + z_2 s_2$. Now, based on the additive property of kurtosis, we have $\text{kurt}(y) = \text{kurt}(z_1 s_1) + \text{kurt}(z_2 s_2) = z_1^4 \text{kurt}(s_1) + z_2^4 \text{kurt}(s_2)$. On the other hand, we made the constraint that the variance of y is equal to 1, based on the same assumption concerning s_1, s_2 . This implies a constraint on \mathbf{z} : $E\{y^2\} = z_1^2 + z_2^2 = 1$. Geometrically, this means that vector \mathbf{z} is constrained to the unit circle on the 2-dimensional plane. The optimization problem is now: what are the maxima of the function $|\text{kurt}(y)| = |z_1^4 \text{kurt}(s_1) + z_2^4 \text{kurt}(s_2)|$ on the unit circle? For simplicity, you may consider that the kurtosis is of the same sign, in which case it absolute value operators can be omitted. The graph of this function is the “optimization landscape” for the problem.

It is not hard to show (Delfosse & Loubaton, 1995) that the maxima are at the points when exactly one of the elements of vector \mathbf{z} is zero and the other non-zero; because of the unit circle constraint, the non-zero element must be equal to 1 or -1 . But these points are exactly the ones when y equals one of the independent components $\pm s_i$, and the problem has been solved.

In practice we would start from some weight vector \mathbf{w} , compute the direction in which the kurtosis of $y = \mathbf{w}^T \mathbf{x}$ is growing most strongly (if kurtosis is positive) or decreasing most strongly (if kurtosis is negative) based on the available sample $\mathbf{x}(1), \dots, \mathbf{x}(T)$ of mixture vector \mathbf{x} , and use a gradient method or one of their extensions for finding a new vector \mathbf{w} . The example can be generalized to arbitrary dimensions, showing that kurtosis can theoretically be used as an optimization criterion for the ICA problem.

However, kurtosis has also some drawbacks in practice, when its value has to be estimated from a measured sample. The main problem is that kurtosis can be very sensitive to outliers (Huber, 1985). Its value may depend on only a few observations in the tails of the distribution, which may be erroneous or irrelevant observations. In other words, kurtosis is not a robust measure of non-Gaussianity.

Thus, other measures of non-Gaussianity might be better than kurtosis in some situations. Below we shall consider negentropy whose properties are rather opposite to those of kurtosis, and finally introduce approximations of negentropy that more or less combine the good properties of both measures.

4.2.2. Negentropy

A second very important measure of non-Gaussianity is given by negentropy. Negentropy is based on the information-theoretic quantity of (differential) entropy.

Entropy is the basic concept of information theory. The entropy of a random variable can be interpreted as the

degree of information that the observation of the variable gives. The more “random”, i.e. unpredictable and unstructured the variable is, the larger its entropy. More rigorously, entropy is closely related to the coding length of the random variable, in fact, under some simplifying assumptions, entropy is the coding length of the random variable. For introductions on information theory, see Cover and Thomas (1991) and Papoulis (1991).

Entropy H is defined for a discrete random variable Y as

$$H(Y) = - \sum_i P(Y = a_i) \log P(Y = a_i) \quad (20)$$

where the a_i are the possible values of Y . This very well-known definition can be generalized for continuous-valued random variables and vectors, in which case it is often called differential entropy. The differential entropy H of a random vector \mathbf{y} with density $f(\mathbf{y})$ is defined as (Cover & Thomas, 1991; Papoulis, 1991):

$$H(\mathbf{y}) = - \int f(\mathbf{y}) \log f(\mathbf{y}) d\mathbf{y}. \quad (21)$$

A fundamental result of information theory is that a Gaussian variable has the largest entropy among all random variables of equal variance. For a proof, see e.g. Cover and Thomas (1991) and Papoulis (1991). This means that entropy could be used as a measure of non-Gaussianity. In fact, this shows that the Gaussian distribution is the “most random” or the least structured of all distributions. Entropy is small for distributions that are clearly concentrated on certain values, i.e. when the variable is clearly clustered, or has a pdf that is very “spiky”.

To obtain a measure of non-Gaussianity that is zero for a Gaussian variable and always non-negative, one often uses a slightly modified version of the definition of differential entropy, called negentropy. Negentropy J is defined as follows:

$$J(\mathbf{y}) = H(\mathbf{y}_{\text{gauss}}) - H(\mathbf{y}) \quad (22)$$

where $\mathbf{y}_{\text{gauss}}$ is a Gaussian random variable of the same covariance matrix as \mathbf{y} . Due to the above-mentioned properties, negentropy is always non-negative, and it is zero if and only if \mathbf{y} has a Gaussian distribution. Negentropy has the additional interesting property that it is invariant for invertible linear transformations (Comon, 1994; Hyvärinen, 1999e).

The advantage of using negentropy, or, equivalently, differential entropy, as a measure of non-Gaussianity is that it is well justified by statistical theory. In fact, negentropy is in some sense the optimal estimator of non-Gaussianity, as far as statistical properties are concerned. The problem in using negentropy is, however, that it is computationally very difficult. Estimating negentropy using the definition would require an estimate (possibly non-parametric) of the pdf. Therefore, simpler approximations of negentropy are very useful, as will be discussed next.

4.2.3. Approximations of negentropy

The estimation of negentropy is difficult, as mentioned above, and therefore this contrast function remains mainly a theoretical one. In practice, some approximations have to be used. Here we introduce approximations that have very promising properties, and which will be used in the following to derive an efficient method for ICA.

The classical method of approximating negentropy is using higher-order moments, for example as follows (Jones & Sibson, 1987):

$$J(y) \approx \frac{1}{12} E\{y^3\}^2 + \frac{1}{48} \text{kurt}(y)^2 \quad (23)$$

The random variable y is assumed to be of zero mean and unit variance. However, the validity of such approximations may be rather limited. In particular, these approximations suffer from the non-robustness encountered with kurtosis.

To avoid the problems encountered with the preceding approximations of negentropy, new approximations were developed in Hyvärinen (1998b). These approximations were based on the maximum-entropy principle. In general we obtain the following approximation:

$$J(y) \approx \sum_{i=1}^p k_i [E\{G_i(y)\} - E\{G_i(\nu)\}]^2, \quad (24)$$

where k_i are some positive constants, and ν is a Gaussian variable of zero mean and unit variance (i.e. standardized). The variable y is assumed to be of zero mean and unit variance, and the functions G_i are some non-quadratic functions (Hyvärinen, 1998b). Note that even in cases where this approximation is not very accurate, Eq. (24) can be used to construct a measure of non-Gaussianity that is consistent in the sense that it is always non-negative, and equal to zero if y has a Gaussian distribution.

In the case where we use only one non-quadratic function G , the approximation becomes

$$J(y) \propto [E\{G(y)\} - E\{G(\nu)\}]^2 \quad (25)$$

for practically any non-quadratic function G . This is clearly a generalization of the moment-based approximation in Eq. (23), if y is symmetric. Indeed, taking $G(y) = y^4$, one then obtains exactly Eq. (23), i.e. a kurtosis-based approximation.

But the point here is that by choosing G wisely, one obtains approximations of negentropy that are much better than the one given by Eq. (23). In particular, choosing G that does not grow too fast, one obtains more robust estimators. The following choices of G have proved very useful:

$$G_1(u) = \frac{1}{a_1} \log \cosh a_1 u, \quad G_2(u) = -\exp(-u^2/2) \quad (26)$$

where $1 \leq a_1 \leq 2$ is some suitable constant.

Thus we obtain approximations of negentropy that give a very good compromise between the properties of the two classical non-Gaussianity measures given by kurtosis and negentropy. They are conceptually simple, fast to compute,

yet have appealing statistical properties, especially robustness. Therefore, we shall use these contrast functions in our ICA methods. Since kurtosis can be expressed in this same framework, it can still be used by our ICA methods. A practical algorithm based on these contrast functions will be presented in Section 6.

4.3. Minimization of mutual information

Another approach for ICA estimation, inspired by information theory, is minimization of mutual information. We will explain this approach here, and show that it leads to the same principle of finding most non-Gaussian directions as was described above. In particular, this approach gives a rigorous justification for the heuristic principles used above.

4.3.1. Mutual information

Using the concept of differential entropy, we define the mutual information I between m (scalar) random variables, y_i , $i = 1 \dots m$ as follows:

$$I(y_1, y_2, \dots, y_m) = \sum_{i=1}^m H(y_i) - H(y). \quad (27)$$

Mutual information is a natural measure of the dependence between random variables. In fact, it is equivalent to the well-known Kullback–Leibler divergence between the joint density $f(\mathbf{y})$ and the product of its marginal densities; a very natural measure for independence. It is always non-negative, and zero if and only if the variables are statistically independent. Thus, mutual information takes into account the whole dependence structure of the variables, and not only the covariance, like PCA and related methods.

Mutual information can be interpreted by using the interpretation of entropy as code length. The terms $H(y_i)$ give the lengths of codes for the y_i when these are coded separately, and $H(\mathbf{y})$ gives the code length when \mathbf{y} is coded as a random vector, i.e. all the components are coded in the same code. Mutual information thus shows what code length reduction is obtained by coding the whole vector instead of the separate components. In general, better codes can be obtained by coding the whole vector. However, if the y_i are independent, they give no information on each other, and one could just as well code the variables separately without increasing code length.

An important property of mutual information (Cover & Thomas, 1991; Papoulis, 1991) is that we have for an invertible linear transformation $\mathbf{y} = \mathbf{W}\mathbf{x}$

$$I(y_1, y_2, \dots, y_n) = \sum_i H(y_i) - H(\mathbf{x}) - \log|\det \mathbf{W}|. \quad (28)$$

Now, let us consider what happens if we constrain the y_i to be *uncorrelated* and of unit variance. This means

$E\{\mathbf{y}\mathbf{y}^T\} = \mathbf{W}E\{\mathbf{x}\mathbf{x}^T\}\mathbf{W}^T = \mathbf{I}$, which implies

$$\det \mathbf{I} = 1 = (\det \mathbf{W}E\{\mathbf{x}\mathbf{x}^T\}\mathbf{W}^T) \\ = (\det \mathbf{W})(\det E\{\mathbf{x}\mathbf{x}^T\})(\det \mathbf{W}^T), \quad (29)$$

and this implies that $\det \mathbf{W}$ must be constant. Moreover, for y_i of unit variance, entropy and negentropy differ only by a constant, and the sign. Thus we obtain,

$$I(y_1, y_2, \dots, y_n) = C - \sum_i J(y_i). \quad (30)$$

where C is a constant that does not depend on \mathbf{W} . This shows the fundamental relation between negentropy and mutual information.

4.3.2. Defining ICA by mutual information

Since mutual information is the natural information-theoretic measure of the independence of random variables, we could use it as the criterion for finding the ICA transform. In this approach that is an alternative to the model estimation approach, we define the ICA of a random vector \mathbf{x} as an invertible transformation as in Eq. (6), where the matrix \mathbf{W} is determined so that the mutual information of the transformed components s_i is minimized.

It is now obvious from Eq. (30) that finding an invertible transformation \mathbf{W} that minimizes the mutual information is roughly equivalent to finding directions in which the negentropy is maximized. More precisely, it is roughly equivalent to finding 1-D subspaces such that the projections in those subspaces have maximum negentropy. Rigorously, speaking, Eq. (30) shows that ICA estimation by minimization of mutual information is equivalent to maximizing the sum of non-Gaussianities of the estimates, when the estimates are constrained to be uncorrelated. The constraint of uncorrelatedness is in fact not necessary, but simplifies the computations considerably, as one can then use the simpler form in Eq. (30) instead of the more complicated form in Eq. (28).

Thus, we see that the formulation of ICA as minimization of mutual information gives another rigorous justification of our more heuristically introduced idea of finding maximally non-Gaussian directions.

4.4. Maximum likelihood estimation

4.4.1. The likelihood

A very popular approach for estimating the ICA model is maximum likelihood estimation, which is closely connected to the infomax principle. Here we discuss this approach, and show that it is essentially equivalent to minimization of mutual information.

It is possible to formulate directly the likelihood in the noise-free ICA model, which was done in Pham, Garrat and Jutten (1992), and then estimate the model by a maximum likelihood method. Denoting by $\mathbf{W} = (\mathbf{w}_1, \dots, \mathbf{w}_n)^T$ the matrix \mathbf{A}^{-1} , the log-likelihood takes the form (Pham et al.,

1992):

$$L = \sum_{t=1}^T \sum_{i=1}^n \log f_i(\mathbf{w}_i^T \mathbf{x}(t)) + T \log |\det \mathbf{W}| \quad (31)$$

where the f_i are the density functions of the s_i (here assumed to be known), and the $\mathbf{x}(t)$, $t = 1, \dots, T$ are the realizations of \mathbf{x} . The term $\log |\det \mathbf{W}|$ in the likelihood comes from the classic rule for (linearly) transforming random variables and their densities (Papoulis, 1991). In general, for any random vector \mathbf{x} with density p_x and for any matrix \mathbf{W} , the density of $\mathbf{y} = \mathbf{W}\mathbf{x}$ is given by $p_x(\mathbf{W}^{-1}\mathbf{y})|\det \mathbf{W}^{-1}|$.

4.4.2. The infomax principle

Another related contrast function was derived from a neural network viewpoint in Bell and Sejnowski (1995) and Nadal and Parga (1994). This was based on maximizing the output entropy (or information flow) of a neural network with non-linear outputs. Assume that \mathbf{x} is the input to the neural network whose outputs are of the form $\phi_i(\mathbf{w}_i^T \mathbf{x})$, where the ϕ_i are some non-linear scalar functions, and the \mathbf{w}_i are the weight vectors of the neurons. One then wants to maximize the entropy of the outputs:

$$L_2 = H(\phi_1(\mathbf{w}_1^T \mathbf{x}), \dots, \phi_n(\mathbf{w}_n^T \mathbf{x})). \quad (32)$$

If the ϕ_i are well chosen, this framework also enables the estimation of the ICA model. Indeed, several authors, e.g. Cardoso (1997) and Pearlmutter and Parra (1997), proved the surprising result that the principle of network entropy maximization, or “infomax”, is equivalent to maximum likelihood estimation. This equivalence requires that the non-linearities ϕ_i used in the neural network are chosen as the cumulative distribution functions corresponding to the densities f_i , i.e. $\phi_i'(\cdot) = f_i(\cdot)$.

4.4.3. Connection to mutual information

To see the connection between likelihood and mutual information, consider the expectation of the log-likelihood:

$$\frac{1}{T} E\{L\} = \sum_{i=1}^n E\{\log f_i(\mathbf{w}_i^T \mathbf{x})\} + \log |\det \mathbf{W}|. \quad (33)$$

Actually, if the f_i were equal to the actual distributions of $\mathbf{w}_i^T \mathbf{x}$, the first term would be equal to $-\sum_i H(\mathbf{w}_i^T \mathbf{x})$. Thus the likelihood would be equal, up to an additive constant, to the negative of mutual information as given in Eq. (28).

Actually, in practice the connection is even stronger. This is because in practice we do not know the distributions of the independent components. A reasonable approach would be to estimate the density of $\mathbf{w}_i^T \mathbf{x}$ as part of the ML estimation method, and use this as an approximation of the density of s_i . In this case, likelihood and mutual information are, for all practical purposes, equivalent.

Nevertheless, there is a small difference that may be very important in practice. The problem with maximum likelihood estimation is that the densities f_i must be estimated correctly. They need not be estimated with any great

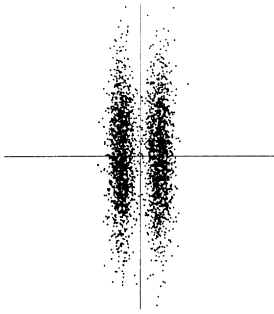


Fig. 9. An illustration of projection pursuit and the “interestingness” of non-Gaussian projections. The data in this figure is clearly divided into two clusters. However, the principal component, i.e. the direction of maximum variance, would be vertical, providing no separation between the clusters. In contrast, the strongly non-Gaussian projection pursuit direction is horizontal, providing optimal separation of the clusters.

precision: in fact it is enough to estimate whether they are sub- or super- Gaussian (Cardoso & Laheld, 1996; Hyvärinen & Oja, 1998; Lee, Girolami & Sejnowski, 1999). In many cases, in fact, we have enough prior knowledge on the independent components, and we do not need to estimate their nature from the data. In any case, if the information on the nature of the independent components is not correct, ML estimation will give completely wrong results. Some care must be taken with ML estimation, therefore. In contrast, using reasonable measures of non-Gaussianity, this problem does not usually arise.

4.5. ICA and projection pursuit

It is interesting to note how our approach to ICA makes explicit the connection between ICA and projection pursuit. Projection pursuit (Friedman, 1987; Friedman & Tukey, 1974; Huber, 1985; Jones & Sibson, 1987) is a technique developed in statistics for finding “interesting” projections of multidimensional data. Such projections can then be used for optimal visualization of the data, and for such purposes as density estimation and regression. In basic (1-D) projection pursuit, we try to find directions such that the projections of the data in those directions have interesting distributions, i.e. display some structure. It has been argued by Huber (1985) and Jones and Sibson (1987) that the Gaussian distribution is the least interesting one, and that the most interesting directions are those that show the least Gaussian distribution. This is exactly what we do to estimate the ICA model.

The usefulness of finding such projections can be seen in Fig. 9, where the projection on the projection pursuit direction, which is horizontal, clearly shows the clustered structure of the data. The projection on the first principal component (vertical), on the other hand, fails to show this structure.

Thus, in the general formulation, ICA can be considered a variant of projection pursuit. All the non-Gaussianity measures and the corresponding ICA algorithms presented

here could also be called projection pursuit “indices” and algorithms. In particular, the projection pursuit allows us to tackle the situation where there are less independent components s_i than original variables x_i . Assuming that those dimensions of the space that are not spanned by the independent components are filled by Gaussian noise, we see that computing the non-Gaussian projection pursuit directions, we effectively estimate the independent components. When all the non-Gaussian directions have been found, all the independent components have been estimated. Such a procedure can be interpreted as a hybrid of projection pursuit and ICA.

However, it should be noted that in the formulation of projection pursuit, no data model or assumption about independent components is made. If the ICA model holds, optimizing the ICA non-Gaussianity measures produce independent components; if the model does not hold, then what we get are the projection pursuit directions.

5. Preprocessing for ICA

In the preceding section, we discussed the statistical principles underlying ICA methods. Practical algorithms based on these principles will be discussed in the next section. However, before applying an ICA algorithm on the data, it is usually very useful to do some preprocessing. In this section, we discuss some preprocessing techniques that make the problem of ICA estimation simpler and better conditioned.

5.1. Centering

The most basic and necessary preprocessing is to center \mathbf{x} , i.e. subtract its mean vector $\mathbf{m} = E\{\mathbf{x}\}$ so as to make \mathbf{x} a zero-mean variable. This implies that \mathbf{s} is zero-mean as well, as can be seen by taking expectations on both sides of Eq. (4).

This preprocessing is made solely to simplify the ICA algorithms: It does not mean that the mean could not be estimated. After estimating the mixing matrix \mathbf{A} with centered data, we can complete the estimation by adding the mean vector of \mathbf{s} back to the centered estimates of \mathbf{s} . The mean vector is given by $\mathbf{A}^{-1}\mathbf{m}$, where \mathbf{m} is the mean that was subtracted in the preprocessing.

5.2. Whitening

Another useful preprocessing strategy in ICA is to first whiten the observed variables. This means that before the application of the ICA algorithm (and after centering), we transform the observed vector \mathbf{x} linearly so that we obtain a new vector $\tilde{\mathbf{x}}$ which is white, i.e. its components are uncorrelated and their variances equal unity. In other words, the covariance matrix of $\tilde{\mathbf{x}}$ equals the identity matrix:

$$E\{\tilde{\mathbf{x}}\tilde{\mathbf{x}}^T\} = \mathbf{I}. \quad (34)$$

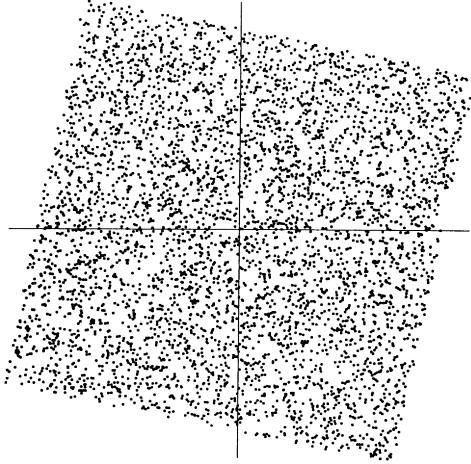


Fig. 10. The joint distribution of the whitened mixtures.

The whitening transformation is always possible. One popular method for whitening is to use the eigenvalue decomposition (EVD) of the covariance matrix $E\{\mathbf{x}\mathbf{x}^T\} = \mathbf{E}\mathbf{D}\mathbf{E}^T$, where \mathbf{E} is the orthogonal matrix of eigenvectors of $E\{\mathbf{x}\mathbf{x}^T\}$ and \mathbf{D} is the diagonal matrix of its eigenvalues, $\mathbf{D} = \text{diag}(d_1, \dots, d_n)$. Note that $E\{\mathbf{x}\mathbf{x}^T\}$ can be estimated in a standard way from the available sample $\mathbf{x}(1), \dots, \mathbf{x}(T)$. Whitening can now be done by

$$\tilde{\mathbf{x}} = \mathbf{E}\mathbf{D}^{-1/2}\mathbf{E}^T\mathbf{x} \quad (35)$$

where the matrix $\mathbf{D}^{-1/2}$ is computed by a simple component-wise operation as $\mathbf{D}^{-1/2} = \text{diag}(d_1^{-1/2}, \dots, d_n^{-1/2})$. It is easy to check that now $E\{\tilde{\mathbf{x}}\tilde{\mathbf{x}}^T\} = \mathbf{I}$.

Whitening transforms the mixing matrix into a new one, $\tilde{\mathbf{A}}$. We have from Eqs. (4) and (35):

$$\tilde{\mathbf{x}} = \mathbf{E}\mathbf{D}^{-1/2}\mathbf{E}^T\mathbf{A}\mathbf{s} = \tilde{\mathbf{A}}\mathbf{s} \quad (36)$$

The utility of whitening resides in the fact that the new mixing matrix $\tilde{\mathbf{A}}$ is orthogonal. This can be seen from

$$E\{\tilde{\mathbf{x}}\tilde{\mathbf{x}}^T\} = \tilde{\mathbf{A}}E\{\mathbf{s}\mathbf{s}^T\}\tilde{\mathbf{A}}^T = \tilde{\mathbf{A}}\tilde{\mathbf{A}}^T = \mathbf{I}. \quad (37)$$

Here we see that whitening reduces the number of parameters to be estimated. Instead of having to estimate the n^2 parameters that are the elements of the original matrix \mathbf{A} , we only need to estimate the new, orthogonal mixing matrix $\tilde{\mathbf{A}}$. An orthogonal matrix contains $n(n-1)/2$ degrees of freedom. For example, in two dimensions, an orthogonal transformation is determined by a single angle parameter. In larger dimensions, an orthogonal matrix contains only about half of the number of parameters of an arbitrary matrix. Thus one can say that whitening solves half of the problem of ICA. Because whitening is a very simple and standard procedure, much simpler than any ICA algorithms, it is a good idea to reduce the complexity of the problem this way.

It may also be quite useful to reduce the dimension of the data at the same time as we do the whitening. Then we look at the eigenvalues d_j of $E\{\mathbf{x}\mathbf{x}^T\}$ and discard those that are too

small, as is often done in the statistical technique of principal component analysis. This has often the effect of reducing noise. Moreover, dimension reduction prevents overlearning, which can sometimes be observed in ICA (Hyvärinen, Särelä & Vigário, 1999).

A graphical illustration of the effect of whitening can be seen in Fig. 10, in which the data in Fig. 6 has been whitened. The square defining the distribution is now clearly a rotated version of the original square in Fig. 5. All that is left is the estimation of a single angle that gives the rotation.

In the rest of this paper, we assume that the data has been preprocessed by centering and whitening. For simplicity of notation, we denote the preprocessed data just by \mathbf{x} , and the transformed mixing matrix by \mathbf{A} , omitting the tildes.

5.3. Further preprocessing

The success of ICA for a given data set may depend crucially on performing some application-dependent preprocessing steps. For example, if the data consists of time-signals, some band-pass filtering may be very useful. Note that if we filter linearly the observed signals $x_i(t)$ to obtain new signals, say $x_i^*(t)$, the ICA model still holds for $x_i^*(t)$, with the same mixing matrix.

This can be seen as follows. Denote by \mathbf{X} the matrix that contains the observations $\mathbf{x}(1), \dots, \mathbf{x}(T)$ as its columns, and similarly for \mathbf{S} . Then the ICA model can be expressed as:

$$\mathbf{X} = \mathbf{A}\mathbf{S} \quad (38)$$

Now, time filtering of \mathbf{X} corresponds to multiplying \mathbf{X} from the right by a matrix, let us call it \mathbf{M} . This gives

$$\mathbf{X}^* = \mathbf{X}\mathbf{M} = \mathbf{A}\mathbf{S}\mathbf{M} = \mathbf{A}\mathbf{S}^*, \quad (39)$$

which shows that the ICA model remains still valid.

6. The FastICA algorithm

In the preceding sections, we introduced different measures of non-Gaussianity, i.e. objective functions for ICA estimation. In practice, one also needs an algorithm for maximizing the contrast function, for example the one in Eq. (25). In this section, we introduce a very efficient method of maximization suited for this task. It is here assumed that the data is preprocessed by centering and whitening as discussed in the preceding section.

6.1. FastICA for one unit

To begin with, we shall show the one-unit version of FastICA. By a “unit” we refer to a computational unit, eventually an artificial neuron, having a weight vector \mathbf{w} that the neuron is able to update by a learning rule. The FastICA learning rule finds a direction, i.e. a unit vector \mathbf{w} such that the projection $\mathbf{w}^T\mathbf{x}$ maximizes non-Gaussianity. Non-Gaussianity is here measured by the approximation

of negentropy $J(\mathbf{w}^T \mathbf{x})$ given in Eq. (25). Recall that the variance of $\mathbf{w}^T \mathbf{x}$ must here be constrained to unity; for whitened data this is equivalent to constraining the norm of \mathbf{w} to be unity.

The FastICA is based on a fixed-point iteration scheme for finding a maximum of the non-Gaussianity of $\mathbf{w}^T \mathbf{x}$, as measured in Eq. (25) (Hyvärinen, 1999a; Hyvärinen & Oja, 1997). It can be also derived as an approximative Newton iteration (Hyvärinen, 1999a). Denote by g the derivative of the non-quadratic function G used in Eq. (25); for example the derivatives of the functions in Eq. (26) are:

$$g_1(u) = \tanh(a_1 u), \quad g_2(u) = u \exp(-u^2/2) \quad (40)$$

where $1 \leq a_1 \leq 2$ is some suitable constant, often taken as $a_1 = 1$. The basic form of the FastICA algorithm is as follows:

1. Choose an initial (e.g. random) weight vector \mathbf{w} .
2. Let $\mathbf{w}^+ = E\{\mathbf{x}g(\mathbf{w}^T \mathbf{x})\} - E\{g'(\mathbf{w}^T \mathbf{x})\}\mathbf{w}$
3. Let $\mathbf{w} = \mathbf{w}^+ / \|\mathbf{w}^+\|$
4. If not converged, go back to 2.

Note that convergence means that the old and new values of \mathbf{w} point in the same direction, i.e. their dot-product is (almost) equal to 1. It is not necessary that the vector converges to a single point, since \mathbf{w} and $-\mathbf{w}$ define the same direction. This is again because the independent components can be defined only up to a multiplicative sign. Note also that it is here assumed that the data is prewhitened.

The derivation of FastICA is as follows. First note that the maxima of the approximation of the negentropy of $\mathbf{w}^T \mathbf{x}$ are obtained at certain optima of $E\{G(\mathbf{w}^T \mathbf{x})\}$. According to the Kuhn–Tucker conditions (Luenberger, 1969), the optima of $E\{G(\mathbf{w}^T \mathbf{x})\}$ under the constraint $E\{(\mathbf{w}^T \mathbf{x})^2\} = \|\mathbf{w}\|^2 = 1$ are obtained at points where

$$E\{\mathbf{x}g(\mathbf{w}^T \mathbf{x})\} - \beta \mathbf{w} = 0 \quad (41)$$

Let us try to solve this equation by Newton's method. Denoting the function on the left-hand side of Eq. (41) by F , we obtain its Jacobian matrix $JF(\mathbf{w})$ as

$$JF(\mathbf{w}) = E\{\mathbf{x}\mathbf{x}^T g'(\mathbf{w}^T \mathbf{x})\} - \beta \mathbf{I} \quad (42)$$

To simplify the inversion of this matrix, we decide to approximate the first term in Eq. (42). Since the data is sphered, a reasonable approximation seems to be $E\{\mathbf{x}\mathbf{x}^T g'(\mathbf{w}^T \mathbf{x})\} \approx E\{\mathbf{x}\mathbf{x}^T\}E\{g'(\mathbf{w}^T \mathbf{x})\} = E\{g'(\mathbf{w}^T \mathbf{x})\}\mathbf{I}$.

Thus the Jacobian matrix becomes diagonal, and can easily be inverted. Thus we obtain the following approximative Newton iteration:

$$\mathbf{w}^+ = \mathbf{w} - [E\{\mathbf{x}g(\mathbf{w}^T \mathbf{x})\} - \beta \mathbf{w}] / [E\{g'(\mathbf{w}^T \mathbf{x})\} - \beta] \quad (43)$$

This algorithm can be further simplified by multiplying both sides of Eq. (43) by $\beta - E\{g'(\mathbf{w}^T \mathbf{x})\}$. This gives, after algebraic simplification, the FastICA iteration.

In practice, the expectations in FastICA must be replaced

by their estimates. The natural estimates are of course the corresponding sample means. Ideally, all the data available should be used, but this is often not a good idea because the computations may become too demanding. Then the averages can be estimated using a smaller sample, whose size may have a considerable effect on the accuracy of the final estimates. The sample points should be chosen separately at every iteration. If the convergence is not satisfactory, one may then increase the sample size.

6.2. FastICA for several units

The one-unit algorithm of Section 6.1 estimates just one of the independent components, or one projection pursuit direction. To estimate several independent components, we need to run the one-unit FastICA algorithm using several units (e.g. neurons) with weight vectors $\mathbf{w}_1, \dots, \mathbf{w}_n$.

To prevent different vectors from converging to the same maxima we must *decorrelate* the outputs $\mathbf{w}_1^T \mathbf{x}, \dots, \mathbf{w}_n^T \mathbf{x}$ after every iteration. We present here three methods for achieving this. Note that for whitened \mathbf{x} such a decorrelation is equivalent to orthogonalization.

A simple way of achieving decorrelation is a deflation scheme based on a Gram–Schmidt-like decorrelation. This means that we estimate the independent components one by one. When we have estimated p independent components, or p vectors $\mathbf{w}_1, \dots, \mathbf{w}_p$, we run the one-unit fixed-point algorithm for \mathbf{w}_{p+1} , and after every iteration step subtract from \mathbf{w}_{p+1} the “projections” $\mathbf{w}_{p+1}^T \mathbf{w}_j \mathbf{w}_j$, $j = 1, \dots, p$ of the previously estimated p vectors, and then renormalize \mathbf{w}_{p+1} :

$$1. \text{ Let } \mathbf{w}_{p+1} = \mathbf{w}_{p+1} - \sum_{j=1}^p \mathbf{w}_{p+1}^T \mathbf{w}_j \mathbf{w}_j \quad (44)$$

$$2. \text{ Let } \mathbf{w}_{p+1} = \mathbf{w}_{p+1} / \sqrt{\mathbf{w}_{p+1}^T \mathbf{w}_{p+1}}$$

In certain applications, however, it may be desired to use a symmetric decorrelation, in which no vectors are “privileged” over others (Karhunen, Oja, Wang, Vigário & Joutsensalo, 1997). This can be accomplished, e.g. by the classical method involving matrix square roots,

$$\text{Let } \mathbf{W} = (\mathbf{W}\mathbf{W}^T)^{-1/2} \mathbf{W}. \quad (45)$$

where \mathbf{W} is the matrix $(\mathbf{w}_1, \dots, \mathbf{w}_n)^T$ of the vectors, and the inverse square root $(\mathbf{W}\mathbf{W}^T)^{-1/2}$ is obtained from the eigenvalue decomposition of $\mathbf{W}\mathbf{W}^T = \mathbf{F}\mathbf{\Lambda}\mathbf{F}^T$ as $(\mathbf{W}\mathbf{W}^T)^{-1/2} = \mathbf{F}\mathbf{\Lambda}^{-1/2}\mathbf{F}^T$. A simpler alternative is the following iterative algorithm (Hyvärinen, 1999a),

$$1. \text{ Let } \mathbf{W} = \mathbf{W} / \sqrt{\|\mathbf{W}\mathbf{W}^T\|} \quad (46)$$

Repeat 2. until convergence :

$$2. \text{ Let } \mathbf{W} = \frac{3}{2} \mathbf{W} - \frac{1}{2} \mathbf{W}\mathbf{W}^T \mathbf{W}$$

The norm in step 1 can be almost any ordinary matrix norm, e.g. the 2-norm or the largest absolute row (or column) sum (but not the Frobenius norm).

6.3. FastICA and maximum likelihood

Finally, we give a version of FastICA that shows explicitly the connection to the well-known infomax or maximum likelihood algorithm introduced in Amari, Cichocki and Yang (1996), Bell and Sejnowski (1995), Cardoso and Laheld (1996) and Cichocki and Unbehauen (1996). If we express FastICA using the intermediate formula in Eq. (43), and write it in matrix form (see Hyvärinen (1999b) for details), we see that FastICA takes the following form:

$$\mathbf{W}^+ = \mathbf{W} + \text{diag}(\alpha_i)[\text{diag}(\beta_i) + E\{g(\mathbf{y})\mathbf{y}^T\}]\mathbf{W}. \quad (47)$$

where $\mathbf{y} = \mathbf{W}\mathbf{x}$, $\beta_i = -E\{y_i g(y_i)\}$, and $\alpha_i = -1/(\beta_i + E\{g'(y_i)\})$. The matrix \mathbf{W} needs to be orthogonalized after every step. In this matrix version, it is natural to orthogonalize \mathbf{W} symmetrically.

The above version of FastICA could be compared with the stochastic gradient method for maximizing likelihood (Amari et al., 1996; Bell & Sejnowski, 1995; Cardoso & Laheld, 1996; Cichocki & Unbehauen, 1996):

$$\mathbf{W}^+ = \mathbf{W} + \mu[\mathbf{I} + g(\mathbf{y})\mathbf{y}^T]\mathbf{W}. \quad (48)$$

where μ is the learning rate, not necessarily constant in time. Here, g is a function of the pdf of the independent components: $g = f'_i/f_i$ where the f_i is the pdf of an independent component.

Comparing Eqs. (47) and (48), we see that FastICA can be considered as a fixed-point algorithm for maximum likelihood estimation of the ICA data model. For details, see Hyvärinen (1999b). In FastICA, convergence speed is optimized by the choice of the matrices $\text{diag}(\alpha_i)$ and $\text{diag}(\beta_i)$. Another advantage of FastICA is that it can estimate both sub and superGaussian independent components, which is in contrast to ordinary ML algorithms, which only work for a given class of distributions (see Section 4.4).

6.4. Properties of the FastICA algorithm

The FastICA algorithm and the underlying contrast functions have a number of desirable properties when compared with existing methods for ICA.

1. The convergence is cubic (or at least quadratic), under the assumption of the ICA data model (for a proof, see Hyvärinen (1999a)). This is in contrast to ordinary ICA algorithms based on (stochastic) gradient descent methods, where the convergence is only linear. This means a very fast convergence, as has been confirmed by simulations and experiments on real data (see Giannakopoulos, Karhunen and Oja (1998)).
2. Contrary to gradient-based algorithms, there are no step size parameters to choose. This means that the algorithm is easy to use.

3. The algorithm finds directly independent components of (practically) any non-Gaussian distribution using any non-linearity g . This is in contrast to many algorithms, where some estimate of the probability distribution function has to be first available, and the non-linearity must be chosen accordingly.
4. The performance of the method can be optimized by choosing a suitable non-linearity g . In particular, one can obtain algorithms that are robust and/or of minimum variance. In fact, the two non-linearities in Eq. (40) have some optimal properties; for details see Hyvärinen (1999a).
5. The independent components can be estimated one by one, which is roughly equivalent to doing projection pursuit. This is useful in exploratory data analysis, and decreases the computational load of the method in cases where only some of the independent components need to be estimated.
6. The FastICA has most of the advantages of neural algorithms: it is parallel, distributed, computationally simple, and requires little memory space. Stochastic gradient methods seem to be preferable only if fast adaptivity in a changing environment is required.

A Matlab™ implementation of the FastICA algorithm is available on the World Wide Web free of charge.¹

7. Applications of ICA

In this section we review some applications of ICA. The most classical application of ICA, the cocktail-party problem, was already explained in Section 1 of this paper.

7.1. Separation of artifacts in MEG data

Magnetoencephalography (MEG) is a non-invasive technique by which the activity or the cortical neurons can be measured with very good temporal resolution and moderate spatial resolution. When using a MEG record, as a research or clinical tool, the investigator may face a problem of extracting the essential features of the neuromagnetic signals in the presence of artifacts. The amplitude of the disturbances may be higher than that of the brain signals and the artifacts may resemble pathological signals in shape.

In Vigário, Jousmäki, Hamäläinen, Hari and Oja (1998a), the authors introduced a new method to separate brain activity from artifacts using ICA. The approach is based on the assumption that the brain activity and the artifacts, e.g. eye movements or blinks, or sensor malfunctions, are anatomically and physiologically separate processes, and this separation is reflected in the statistical independence between the magnetic signals generated by those processes. The approach follows the earlier experiments with EEG signals, reported in Vigário (1997). A related approach is that of Makeig, Bell, Jung and Sejnowski (1996).

The MEG signals were recorded in a magnetically

¹ WWW address: <http://www.cis.hut.fi/projects/ica/fastica/>

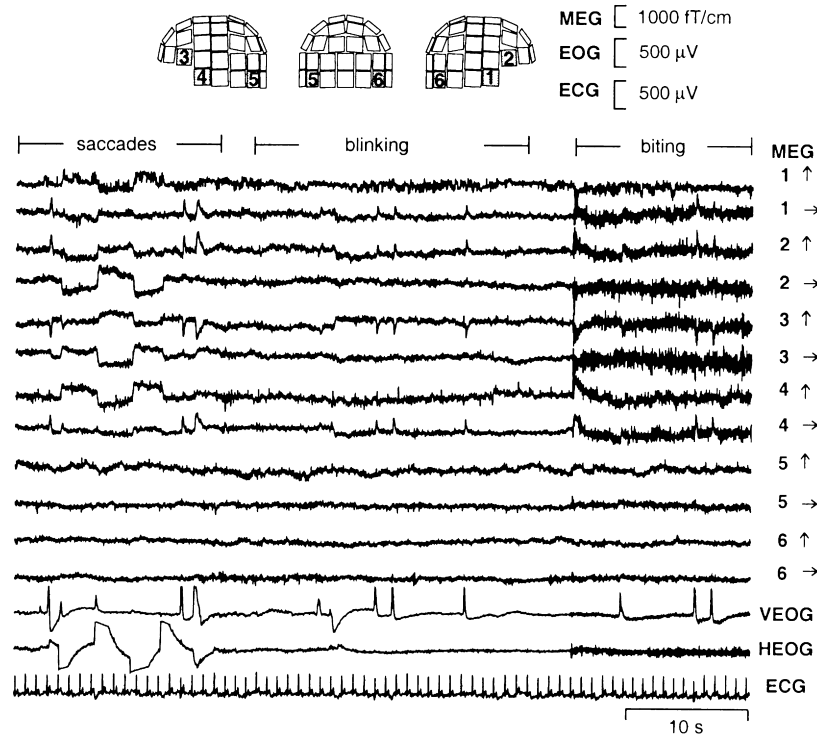


Fig. 11. (From Vigário et al., 1998a.) Samples of MEG signals, showing artifacts produced by blinking, saccades, biting and cardiac cycle. For each of the six positions shown, the two orthogonal directions of the sensors are plotted. Reprinted with permission from the MIT Press.

shielded room with a 122-channel whole-scalp Neuromag-122 neuromagnetometer. This device collects data at 61 locations over the scalp, using orthogonal double-loop pick-up coils that couple strongly to a local source just underneath. The test person was asked to blink and make horizontal saccades, in order to produce typical ocular (eye) artifacts. Moreover, to produce myographic (muscle) artifacts, the subject was asked to bite his teeth for as long as 20 s. Yet another artifact was created by placing a digital watch 1 m away from the helmet into the shielded room.

Fig. 11 presents a subset of 12 spontaneous MEG signals $x_i(t)$ from the frontal, temporal, and occipital areas (Vigário et al., 1998a,b). The figure also shows the positions of the corresponding sensors on the helmet. Due to the dimension of the data (122 magnetic signals were recorded), it is impractical to plot all the MEG signals $x_i(t)$, $i = 1, \dots, 122$. Also two electro-oculogram channels and the electrocardiogram are presented, but they were not used in computing the ICA.

The signal vector \mathbf{x} in the ICA model (4) consists now of the amplitudes $x_i(t)$ of the 122 signals at a certain time point, so the dimensionality is $n = 122$. In the theoretical model, \mathbf{x} is regarded as a random vector, and the measurements $\mathbf{x}(t)$ give a set of realizations of \mathbf{x} as time proceeds. Note that in the basic ICA model that we are using, the temporal correlations in the signals are not utilized at all.

The $\mathbf{x}(t)$ vectors were whitened using PCA and the dimensionality was decreased at the same time. Then, using the FastICA algorithm, a subset of the rows of the separating matrix \mathbf{W} of Eq. (6) were computed. Once a

vector \mathbf{w}_i has become available, an ICA signal $s_i(t)$ can be computed from $s_i(t) = \mathbf{w}_i^T \mathbf{x}(t)$ with $\mathbf{x}(t)$ now denoting the whitened and lower dimensional signal vector.

Fig. 12 shows sections of 9 independent components (IC's) $s_i(t)$, $i = 1, \dots, 9$ found from the recorded data together with the corresponding field patterns (Vigário et al., 1998a). The first two IC's are clearly due to the muscular activity originated from the biting. Their separation into two components seems to correspond, on the basis of the field patterns, to two different sets of muscles that were activated during the process. IC3 and IC5 are showing the horizontal eye movements and the eye blinks, respectively. IC4 represents the cardiac artifact that is very clearly extracted.

To find the remaining artifacts, the data were high-pass filtered, with cutoff frequency at 1 Hz. Next, the independent component IC8 was found. It shows the artifact originated at the digital watch clearly, located to the right side of the magnetometer. The last independent component IC9 is related to a sensor presenting higher RMS (root mean squared) noise than the others.

The results of Fig. 12 clearly show that using the ICA technique and the FastICA algorithm, it is possible to isolate both eye movement and eye blinking artifacts, as well as cardiac, myographic, and other artifacts from MEG signals. The FastICA algorithm is an especially suitable tool, because artifact removal is an interactive technique and the investigator may freely choose how many of the IC's he or she wants.

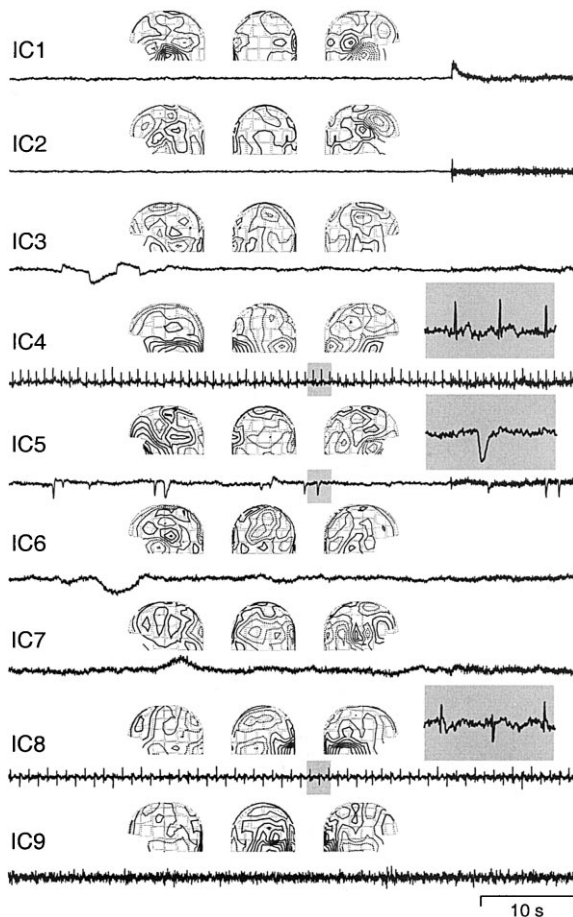


Fig. 12. (From Vigário et al., 1998a.) Nine independent components found from the MEG data. For each component the left, back and right views of the field patterns generated by these components are shown—full line stands for magnetic flux coming out from the head, and dotted line the flux inwards. Reprinted by permission from the MIT Press.

In addition to reducing artifacts, ICA can be used to decompose evoked fields (Vigário et al., 1998b), which enables direct access to the underlying brain functioning, which is likely to be of great significance in neuroscientific research.

7.2. Finding hidden factors in financial data

It is a tempting alternative to try ICA on financial data. There are many situations in that application domain in which parallel time series are available, such as currency exchange rates or daily returns of stocks, that may have some common underlying factors. ICA might reveal some driving mechanisms that otherwise remain hidden. In a recent study of a stock portfolio (Back & Weigend, 1998), it was found that ICA is a complementary tool to PCA, allowing the underlying structure of the data to be more readily observed.

In Kiviluoto and Oja (1998) we applied ICA on a different problem: the cashflow of several stores belonging to the same retail chain, trying to find the fundamental factors

common to all stores that affect the cashflow data. Thus, the cashflow effect of the factors specific to any particular store, i.e. the effect of the actions taken at the individual stores and in its local environment could be analyzed.

The assumption of having some underlying independent components in this specific application may not be unrealistic. For example, factors like seasonal variations due to holidays and annual variations, and factors having a sudden effect on the purchasing power of the customers like prize changes of various commodities, can be expected to have an effect on all the retail stores, and such factors can be assumed to be roughly independent of each other. Yet, depending on the policy and skills of the individual manager like e.g. advertising efforts, the effect of the factors on the cash flow of specific retail outlets are slightly different. By ICA, it is possible to isolate both the underlying factors and the effect weights, thus also making it possible to group the stores on the basis of their managerial policies using only the cash flow time series data.

The data consisted of the weekly cash flow in 40 stores that belong to the same retail chain; the cash flow measurements cover 140 weeks. Some examples of the original data $x_i(t)$ are shown in Fig. 13.

The prewhitening was performed so that the original signal vectors were projected to the subspace spanned by their first five principal components and the variances were normalized to 1. Thus the dimension of the signal space was decreased from 40 to 5. Using the FastICA algorithm, five IC's $s_i(t)$, $i = 1, \dots, 5$ were estimated. As depicted in Fig. 14, the FastICA algorithm has found several clearly different fundamental factors hidden in the original data.

The factors have clearly different interpretations. The upmost two factors follow the sudden changes that are caused by holidays, etc.; the most prominent example is the Christmas time. The factor on the bottom row, on the other hand, reflects the slower seasonal variation, with the effect of the summer holidays clearly visible. The factor on the third row could represent a still slower variation, something resembling a trend. The last factor, on the fourth row, is different from the others; it might be that this factor follows mostly the relative competitive position of the retail chain with respect to its competitors, but other interpretations are also possible.

More details on the experiments and their interpretation can be found in Kiviluoto and Oja (1998).

7.3. Reducing noise in natural images

The third example deals with finding ICA filters for natural images and, based on the ICA decomposition, removing noise from images corrupted with additive Gaussian noise.

A set of digitized natural images was used. Denote the vector of pixel gray levels in an image window by \mathbf{x} . Note that, contrary to the other two applications in the previous sections, we are not this time considering multivalued time

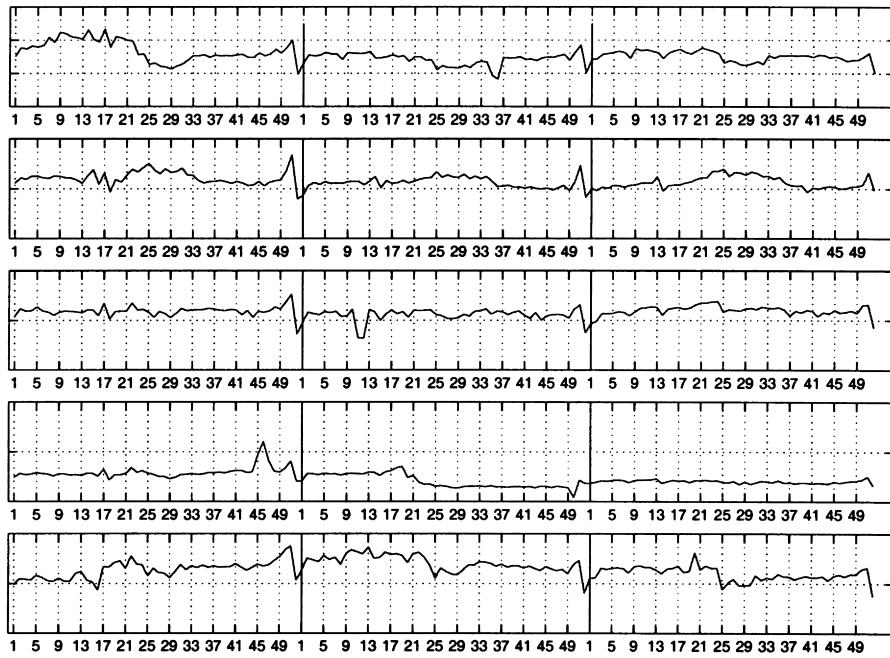


Fig. 13. (Adapted from Kiviluoto and Oja, 1998.) Five samples of the original cashflow time series (mean removed, normalized to unit standard deviation). Horizontal axis: time in weeks.

series or images changing with time; instead the elements of \mathbf{x} are indexed by the location in the image window or patch. The sample windows were taken at random locations. The 2-D structure of the windows is of no significance here: row by row scanning was used to turn a square image window into a vector of pixel values. The independent components of such image windows are represented in Fig. 4. Each window in this figure corresponds to one of the columns

\mathbf{a}_i of the mixing matrix \mathbf{A} . Thus an observed image window is a superposition of these windows as in Eq. (5), with independent coefficients (Bell & Sejnowski, 1997; Olshausen & Field, 1996).

Now, suppose a noisy image model holds:

$$\mathbf{z} = \mathbf{x} + \mathbf{n} \quad (49)$$

where \mathbf{n} is uncorrelated noise, with elements indexed in the

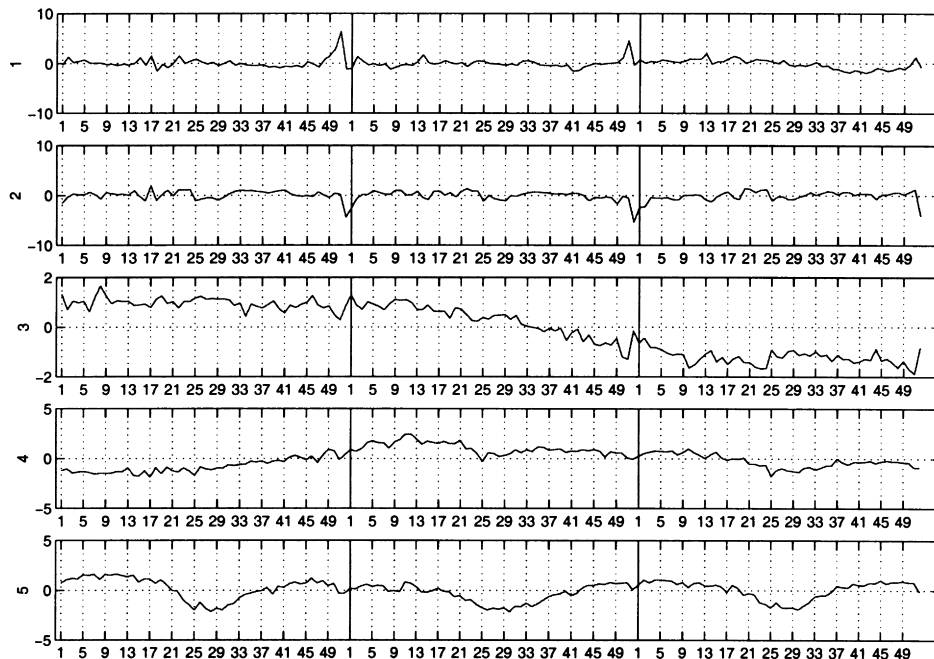


Fig. 14. (Adapted from Kiviluoto and Oja, 1998.) Five independent components or fundamental factors found from the cashflow data.

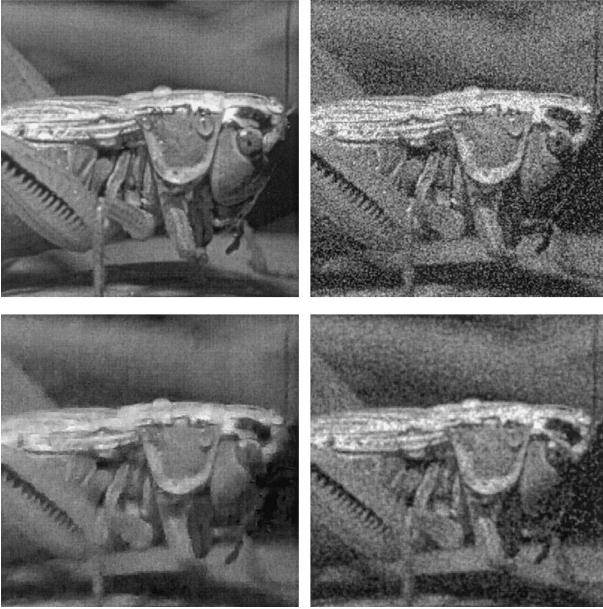


Fig. 15. (From Hyvärinen et al., 2000.) An experiment in denoising. Upper left: original image. Upper right: original image corrupted with noise; the noise level is 50%. Lower left: the recovered image after applying sparse code shrinkage. Lower right: for comparison, a wiener filtered image. Reprinted by permission from the IEEE Press.

image window in the same way as \mathbf{x} , and \mathbf{z} is the measured image window corrupted with noise. Let us further assume that \mathbf{n} is Gaussian and \mathbf{x} is non-Gaussian. There are many ways to clean the noise; one example is to make a transformation to spatial frequency space by DFT, do low-pass filtering, and return to the image space by IDFT (Gonzalez & Wintz, 1987). This is not very efficient, however. A better method is the recently introduced Wavelet Shrinkage method (Donoho, Johnstone & Picard, 1995) in which a transform based on wavelets is used, or methods based on median filtering (Gonzalez & Wintz, 1987). None of these methods is explicitly taking advantage of the image statistics, however.

We have recently introduced another, statistically principled method called Sparse Code Shrinkage (Hyvärinen, 1999d). It is very closely related to independent component analysis. Briefly, if we model the density of \mathbf{x} by ICA, and assume \mathbf{n} Gaussian, then the Maximum Likelihood (ML) solution for \mathbf{x} given the measurement \mathbf{z} can be developed in the signal model (49).

The ML solution can be simply computed, albeit approximately, by using a decomposition that is an orthogonalized version of ICA. The transform is given by

$$\mathbf{Wz} = \mathbf{Wx} + \mathbf{Wn} = \mathbf{s} + \mathbf{Wn}, \quad (50)$$

where \mathbf{W} is here an orthogonal matrix that is the best orthogonal approximation of the inverse of the ICA mixing matrix. The noise term \mathbf{Wn} is still Gaussian and white. With a suitably chosen orthogonal transform \mathbf{W} , however, the density of $\mathbf{Wx} = \mathbf{s}$ becomes highly non-Gaussian, e.g.

super-Gaussian with a high positive kurtosis. This depends of course on the original \mathbf{x} signals, as we are assuming in fact that there exists a model $\mathbf{x} = \mathbf{W}^T \mathbf{s}$ for the signal, such that the “source signals” or elements of \mathbf{s} have a positive kurtotic density, in which case the ICA transform gives highly superGaussian components. This seems to hold at least for image windows of natural scenes (Mallat, 1989).

It was shown in Hyvärinen (1999d) that, assuming a Laplacian density for s_i , the ML solution for s_i is given by a “shrinkage function” $\hat{s}_i = g([\mathbf{Wz}]_i)$, or in vector form, $\hat{\mathbf{s}} = g(\mathbf{Wz})$. Function $g(\cdot)$ has a characteristic shape: it is zero close to the origin and then linear after a cutting value depending on the parameters of the Laplacian density and the Gaussian noise density. Assuming other forms for the densities, other optimal shrinkage functions can be derived (Hyvärinen, 1999d).

In the Sparse Code Shrinkage method, the shrinkage operation is performed in the rotated space, after which the estimate for the signal in the original space is given by rotating back:

$$\hat{\mathbf{x}} = \mathbf{W}^T \hat{\mathbf{s}} = \mathbf{W}^T g(\mathbf{Wz}). \quad (51)$$

Thus we get the Maximum Likelihood estimate for the image window in which much of the noise has been removed.

The rotation operator \mathbf{W} is such that the sparsity of the components $\mathbf{s} = \mathbf{Wx}$ is maximized. This operator can be learned with a modification of the FastICA algorithm; see Hyvärinen (1999d) for details.

A noise cleaning result is shown in Fig. 15. A noiseless image and a noisy version, in which the noise level is 50% of the signal level, are shown. The results of the Sparse Code Shrinkage method and classic wiener filtering are given, indicating that Sparse Code Shrinkage may be a promising approach. The noise is reduced without blurring edges or other sharp features as much as in wiener filtering. This is largely due to the strongly non-linear nature of the shrinkage operator that is optimally adapted to the inherent statistics of natural images.

7.4. Telecommunications

Finally, we mention another emerging application area of great potential: telecommunications. An example of a real-world communications application where blind separation techniques are useful is the separation of the user’s own signal from the interfering other users’ signals in CDMA (Code-Division Multiple Access) mobile communications (Ristaniemi & Joutsensalo, 1999). This problem is semi-blind in the sense that certain additional prior information is available on the CDMA data model. But the number of parameters to be estimated is often so high that suitable blind source separation techniques taking into account the available prior knowledge provide a clear performance improvement over more traditional estimation techniques (Ristaniemi & Joutsensalo, 1999; Cristescu et al., 2000).

8. Conclusion

ICA is a very general-purpose statistical technique in which observed random data are linearly transformed into components that are maximally independent from each other, and simultaneously have “interesting” distributions. ICA can be formulated as the estimation of a latent variable model. The intuitive notion of maximum non-Gaussianity can be used to derive different objective functions whose optimization enables the estimation of the ICA model. Alternatively, one may use more classical notions like maximum likelihood estimation or minimization of mutual information to estimate ICA; somewhat surprisingly, these approaches are (approximatively) equivalent. A computationally very efficient method performing the actual estimation is given by the FastICA algorithm. Applications of ICA can be found in many different areas such as audio processing, biomedical signal processing, image processing, telecommunications, and econometrics.

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