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Transient Performance Analysis of Zero-Attracting LMS

Jie Chen, Member, IEEE, Cédric Richard, Senior Member, IEEE, Yingying Song, David Brie, Member, IEEE

Abstract—Zero-attracting least-mean-square (ZA-LMS) algorithm has been widely used for online sparse system identification. It combines the LMS framework and ℓ_1 -norm regularization to promote sparsity, and relies on subgradient iterations. Despite the significant interest in ZA-LMS, few works analyzed its transient behavior. The main difficulty lies in the nonlinearity of the update rule. In this work, a detailed analysis in the mean and mean-square sense is carried out in order to examine the behavior of the algorithm. Simulation results illustrate the accuracy of the model and highlight its performance through comparisons with an existing model.

Index Terms—Sparse system identification, zero-attracting LMS, performance analysis, transient behavior

I. INTRODUCTION

The need for finding sparse solutions to system identification problems has been motivated by many applications in signal and image processing [1], [2], such as those in biomedical imaging [3], remote sensing [4] and communications [5]. For instance, channel estimation problems have emphasized the need for sparsity-aware algorithms because channels with a sparse impulse response arise in a number of practical situations. In all cases, the underlying problem can be expressed in terms of finding the ℓ_0 -based sparsest solution, i.e., to find a minimal number of non-zero coefficients that represent the solution of interest. Despite the problem being generally NP-hard, compressive sensing theory has been developed to provides a robust framework to find sparse solutions with low computational complexity. Regularization with ℓ_1 -norm is popularly used as a surrogate of the nonconvex ℓ_0 -norm to promote the sparsity [6].

Built on this principe, zero-attracting LMS (ZA-LMS) was proposed in [7]. The algorithm combines the LMS framework and ℓ_1 -norm regularization, and relies on subgradient iterations. Its stability was analyzed in [7], [8]. Despite the significant interest in ZA-LMS, few works analyzed its behavior. The main difficulty lies in the nonlinearity of the update rule, which arises from the subgradient of the ℓ_1 -norm term. Most related works include [8] and [9]. In [9], the analysis of the transient behavior of ZA-LMS is limited to white input signals. The proposed model involves terms that cannot be explicitly evaluated. In [8], the authors propose a more comprehensive analysis. This work evaluates the mean weight behavior in an exact manner under the Gaussian assumption on estimated weights (or, equivalently, weight errors). Nevertheless, coarse approximations are adopted for the mean-square error analysis. They consist of simply substituting the expectation of the product of two sign functions by the product of their expectations. Finally, another work [10] suggests an analysis of the ℓ_0 regularized LMS by classifying the weights into distinct classes according to their values.

The main difficulty in the analysis of ZA-LMS behavior lies in evaluating expected values that involve the sign of non-zero mean random variables. Making assumptions which are consistent with [7]–[11], we derive an exact model of the transient behavior of ZA-LMS without further approximations. Simulations illustrate the consistency between the simulated results and the theoretical findings, as well as the improved accuracy compared with the previous works.

Notation. Boldface small letters \boldsymbol{x} denote column vectors. Boldface capital letters \boldsymbol{X} denote matrices. $[\boldsymbol{x}]_i$ and $[\boldsymbol{X}]_{ij}$ denote respectively the *i*-th entry of \boldsymbol{x} , and the (i, j)-th entry of \boldsymbol{X} . The superscript $(\cdot)^{\top}$ denotes the transpose of a matrix or a vector. All-zero vector of length N is denoted by $\boldsymbol{0}_N$. The operator sgn $\{\cdot\}$ takes the sign of the entries of its argument. The operator tr $\{\cdot\}$ takes the trace of its matrix argument. The Gaussian distribution with mean μ and variance σ^2 is denoted by $\mathcal{N}(\mu, \boldsymbol{\Sigma})$. The cumulative distribution function (CDF) of the standard Gaussian distribution is denoted by $\phi(\boldsymbol{x})$. The CDF of the multivariate Gaussian distribution with mean μ and covariance $\boldsymbol{\Sigma}$ is denoted by $\Phi(\boldsymbol{x}, \boldsymbol{\mu}, \boldsymbol{\Sigma})$.

II. PROBLEM FORMULATION AND ZA-LMS

Consider an unknown system with input-output relation characterized by the linear model:

$$y_n = \boldsymbol{x}_n^{\top} \boldsymbol{w}^{\star} + z_n \tag{1}$$

with $\boldsymbol{w}^{\star} \in \mathbb{R}^{L}$ an unknown parameter vector, and $\boldsymbol{x}_{n} \in \mathbb{R}^{L}$ the regressor with correlation matrix $\boldsymbol{R}_{x} > 0$ at instant n. The regressor \boldsymbol{x}_{n} and the output signal y_{n} are assumed to be zero-mean and stationary. The modeling error z_{n} is assumed to be stationary, independent and identically distributed (i.i.d.), with zero-mean and variance σ_{z}^{2} , and independent of any other

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J. Chen is with Centre of Intelligent Acoustics and Immersive Communications at School of Marine Science and Technology, Northwestern Polytechinical University, Xi'an, China (email: dr.jie.chen@ieee.org). C. Richard is with the Université Côte d'Azur, OCA, CNRS, France (email: cedric.richard@unice.fr), and on leave at INRIA Sophia Antipolis, France. Y. Song and D. Brie are with Centre de Recherche en Automatique de Nancy (CRAN), Université de Lorraine, CNRS (email: firstname.lastname@univlorraine.fr).

signal. We consider the system identification problem:

$$\boldsymbol{w}^{o} = \arg\min_{\boldsymbol{w}} \mathbb{E}\left\{ [\boldsymbol{y}_{n} - \boldsymbol{w}^{\top} \boldsymbol{x}_{n}]^{2} \right\} + \lambda \|\boldsymbol{w}\|_{1}$$
(2)

where the ℓ_1 -norm is particularly useful for applications where the parameter vector has a sparse structure. In order to solve this regularized problem, zero-attracting LMS was proposed in [7] based on the subgradient iteration:

$$\boldsymbol{w}_{n+1} = \boldsymbol{w}_n + \mu \, \boldsymbol{e}_n \boldsymbol{x}_n - \rho \, \mathrm{sgn}\{\boldsymbol{w}_n\} \tag{3}$$

where $e_n = y_n - \boldsymbol{w}_n^\top \boldsymbol{x}_n$, μ is a positive step size, and $\rho = \mu \lambda$.

III. TRANSIENT BEHAVIOR MODEL OF ZA-LMS

We now study the transient behavior of the ZA-LMS algorithm. Defining the weight error vector \tilde{w}_n as the difference between the estimated weight vector w_n and w^* , namely,

$$\tilde{\boldsymbol{w}}_n = \boldsymbol{w}_n - \boldsymbol{w}^\star, \qquad (4)$$

the analysis of ZA-LMS consists of studying the evolution of the first and second-order moments of \tilde{w}_n over time.

A. Statistical assumptions

We introduce the following statistical assumptions to keep the calculations mathematically tractable:

- A1: The weight-error vector $\tilde{\boldsymbol{w}}_n$ is statistically independent of the input vector \boldsymbol{x}_n .
- **A2:** Any pair of entries $[\tilde{\boldsymbol{w}}_n]_i$ and $[\tilde{\boldsymbol{w}}_n]_j$ with $i \neq j$ is jointly Gaussian.

The independence assumption A1 is widely used in the analysis of adaptive filters [12], [13]. Assumption A2 is consistent with the Gaussian assumptions in [7]–[11]. It is weaker than in [7], [9]–[11] which assume that \tilde{w}_n is Gaussian distributed, but stronger than in [8] which assumes that the entries $[\tilde{w}_n]_i$ are Gaussian random variables. Nevertheless, the authors in [8] use approximations that are ideally suited for independent entries $[\tilde{w}_n]_i$. We illustrate assumption A2 in Fig. 2 with some histograms of $[\tilde{w}_n]_i$ versus $[\tilde{w}_n]_j$. We shall show it is sufficient to make the calculation of the nonlinear terms tractable without further approximations, and to derive a more accurate model.

B. Mean weight behavior model

We first focus on the mean weight behavior of the algorithm. Subtracting \boldsymbol{w}^{\star} from both sides of (3), and using $e_n = z_n - \tilde{\boldsymbol{w}}_n^{\top} \boldsymbol{x}_n$, yields the update relation of $\tilde{\boldsymbol{w}}_n$:

$$\tilde{\boldsymbol{w}}_{n+1} = \tilde{\boldsymbol{w}}_n - \mu \, \boldsymbol{x}_n \boldsymbol{x}_n^\top \tilde{\boldsymbol{w}}_n + \mu \, \boldsymbol{z}_n \boldsymbol{x}_n - \rho \, \mathrm{sgn}\{\boldsymbol{w}^\star + \tilde{\boldsymbol{w}}_n\}$$
(5)

Taking the expectation of (5) and considering A1, we have:

$$\mathbb{E}\{\tilde{\boldsymbol{w}}_{n+1}\} = \mathbb{E}\{\tilde{\boldsymbol{w}}_n\} - \mu \, \boldsymbol{R}_x \mathbb{E}\{\tilde{\boldsymbol{w}}_n\} - \rho \, \mathbb{E}\{\operatorname{sgn}\{\boldsymbol{w}^* + \tilde{\boldsymbol{w}}_n\}\}.$$
(6)

To characterize the evolution of $\mathbb{E}{\{\tilde{w}_{n+1}\}}$ it is necessary to evaluate the last term $\mathbb{E}{\{\operatorname{sgn}\{\boldsymbol{w}^* + \tilde{\boldsymbol{w}}_n\}}$ in (6).

Lemma 1. Consider a random variable $u \sim \mathcal{N}(\mu, \sigma^2)$. The expectation of its sign value is given by:

$$\mathbb{E}\{sgn\{u\}\} = 1 - 2\phi(-\mu/\sigma) \tag{7}$$

As shown in Appendix A, proving this result is almost trivial. Then, the entries of $\mathbb{E}\{\operatorname{sgn}\{w^*+\tilde{w}_n\}\}\$ are obtained by making the following identification:

$$u \leftarrow [\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}_n]_i \tag{8}$$

with

$$\boldsymbol{u} \leftarrow [\boldsymbol{w}^{\star}]_i + \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i\}$$
(9)

$$\sigma^2 \leftarrow \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i^2\} - \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i\}^2 \tag{10}$$

where $\mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i^2\}$ can be extracted from the diagonal of the second-order moment matrix $\boldsymbol{K}_n = \mathbb{E}\{\tilde{\boldsymbol{w}}_n \tilde{\boldsymbol{w}}_n^{\top}\}$ that will be determined in the next subsection.

C. Mean-square error behavior model

Using $e_n = z_n - \tilde{w}_n^{\dagger} x_n$ and considering A1 leads to the following expression for the mean-square error (MSE) [13]:

$$\mathbb{E}\{e_n^2\} \approx \sigma_z^2 + \operatorname{tr}\{\boldsymbol{R}_x \boldsymbol{K}_n\}$$
(11)

where tr{ $R_x K_n$ } is the excess mean-square error (EMSE). We thus need to determine a recursion for K_n in order to evaluate the MSE (or EMSE). Post-multiplying (5) by its transpose, taking the expected value, and using A1 leads to:

$$\boldsymbol{K}_{n+1} = \boldsymbol{K}_n + \mu^2 \sigma_z^2 \boldsymbol{R}_x + \mu^2 \boldsymbol{Q}_1 + \rho^2 \boldsymbol{Q}_2 - \mu (\boldsymbol{Q}_3 + \boldsymbol{Q}_3^\top) - \rho (\boldsymbol{Q}_4 + \boldsymbol{Q}_4^\top) + \mu \rho (\boldsymbol{Q}_5 + \boldsymbol{Q}_5^\top)$$
(12)

where:

$$\boldsymbol{Q}_1 = \mathbb{E}\{\boldsymbol{x}_n \boldsymbol{x}_n^\top \tilde{\boldsymbol{w}}_n \tilde{\boldsymbol{w}}_n^\top \boldsymbol{x}_n \boldsymbol{x}_n^\top\}$$
(13)

$$\boldsymbol{Q}_2 = \mathbb{E}\{\operatorname{sgn}\{\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}_n\}\operatorname{sgn}^{\top}\{\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}_n\}\}$$
(14)

$$\boldsymbol{Q}_3 = \mathbb{E}\{\tilde{\boldsymbol{w}}_n \tilde{\boldsymbol{w}}_n^{\top} \boldsymbol{x}_n \boldsymbol{x}_n^{\top}\}$$
(15)

$$\boldsymbol{Q}_4 = \mathbb{E}\{\tilde{\boldsymbol{w}}_n \operatorname{sgn}^{\top}\{\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}\}\}$$
(16)

$$\boldsymbol{Q}_{5} = \mathbb{E}\{\boldsymbol{x}_{n}\boldsymbol{x}_{n}^{\top}\tilde{\boldsymbol{w}}_{n}\mathrm{sgn}^{\top}\{\boldsymbol{w}^{\star}+\tilde{\boldsymbol{w}}_{n}\}\}.$$
 (17)

We shall now calculate the terms Q_i defined in (13)–(17). Assumption A1 allows to separate x_n and \tilde{w}_n in (13) to help in evaluating Q_1 . No closed-form expression can, however, be provided for if no distribution for input data x_n is assumed. In order to evaluate the fourth-order moments of x_n , and then provide a closed-form expression for Q_1 , we consider the special case where x_n is zero-mean Gaussian. Using A1 and Isserlis' theorem (decomposition of higher-order moments of the multivariate Gaussian distribution), we then obtain:

$$\boldsymbol{Q}_1 = 2\boldsymbol{R}_x \boldsymbol{K}_n \boldsymbol{R}_x + \operatorname{tr} \{ \boldsymbol{R}_x \boldsymbol{K}_n \} \boldsymbol{R}_x$$
(18)

The rest of the analysis, which consists of calculating the terms Q_2 to Q_5 in (14)–(17), remains valid even if the Gaussian assumption on x_n is relaxed. To evaluate Q_2 , we introduce the following lemma.

Lemma 2. Consider two random variables u and v which are jointly Gaussian, namely,

$$\begin{bmatrix} u \\ v \end{bmatrix} \sim \mathcal{N} \left(\boldsymbol{\mu} := \begin{bmatrix} \mu_u \\ \mu_v \end{bmatrix}, \, \boldsymbol{\Sigma}_{uv} := \begin{bmatrix} \sigma_u^2 & \rho_{uv} \\ \rho_{uv} & \sigma_v^2 \end{bmatrix} \right) \tag{19}$$

with μ and Σ_{uv} their mean vector and covariance matrix, defined as described above. The expectation $\mathbb{E}\{sgn\{u\}sgn\{v\}\}\$ is given by:

$$\mathbb{E}\{sgn\{u\}sgn\{v\}\}\$$

$$=\Phi(\mathbf{0}_{2}, [\mu_{u}, \mu_{v}]^{\top}, \mathbf{\Sigma}_{uv}) + \Phi(\mathbf{0}_{2}, -[\mu_{u}, \mu_{v}]^{\top}, \mathbf{\Sigma}_{uv}) \quad (20)$$

$$-\Phi(\mathbf{0}_{2}, [\mu_{u}, -\mu_{v}]^{\top}, \overline{\mathbf{\Sigma}}_{uv}) - \Phi(\mathbf{0}_{2}, [-\mu_{u}, \mu_{v}]^{\top}, \overline{\mathbf{\Sigma}}_{uv})$$

with $\overline{\Sigma}_{uv} = \Sigma_{uv} \circ \begin{bmatrix} 1 & -1 \\ -1 & 1 \end{bmatrix}$, and \circ the element-wise product.



Fig. 2. Histograms of bivariate vector $[[\tilde{\boldsymbol{w}}_n]_i, [\tilde{\boldsymbol{w}}_n]_i]$ with 5000 samples.

See the proof in Appendix B. The diagonal entries of Q_2 are simply $[Q_2]_{ii} = 1$. The entries $[Q_2]_{ij}$ for $i \neq j$ are obtained by making the following identification:

$$u \leftarrow [\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}_n]_i \tag{21}$$

$$v \leftarrow [\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}_n]_j \tag{22}$$

with

$$\mu_u \leftarrow \mathbb{E}\{[\boldsymbol{w}^\star + \tilde{\boldsymbol{w}}_n]_i\}$$
(23)

$$\mu_v \leftarrow \mathbb{E}\{[\boldsymbol{w}^* + \tilde{\boldsymbol{w}}_n]_j\}$$
(24)

$$\sigma_u^2 \leftarrow \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i^2\} - \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i\}^2 \tag{25}$$

$$\sigma_v^2 \leftarrow \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_j^2\} - \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_j\}^2 \tag{26}$$

$$\rho_{uv} \leftarrow \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i[\tilde{\boldsymbol{w}}_n]_j\} - \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i\}\mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_j\}$$
(27)

where $\mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i[\tilde{\boldsymbol{w}}_n]_j\}$ can be extracted from $[\boldsymbol{K}_n]_{ij}$. Considering A1, the term \boldsymbol{Q}_3 in (15) is given by:

$$\boldsymbol{Q}_3 = \boldsymbol{K}_n \boldsymbol{R}_x. \tag{28}$$

In order to calculate Q_4 , we introduce the following lemma.

Lemma 3. Consider two random variables u and v which are jointly Gaussian, namely,

$$\begin{bmatrix} u \\ v \end{bmatrix} \sim \mathcal{N} \left(\boldsymbol{\mu} := \begin{bmatrix} \mu_u \\ \mu_v \end{bmatrix}, \, \boldsymbol{\Sigma}_{uv} := \begin{bmatrix} \sigma_u^2 & \rho_{uv} \\ \rho_{uv} & \sigma_v^2 \end{bmatrix} \right)$$
(29)

with μ and Σ_{uv} their mean vector and covariance matrix, defined as described above. The expectation $\mathbb{E}\{u \operatorname{sgn}\{v\}\}$ is given by the expressions in Appendix C.

It is worth noting that, in Lemma 2 and Lemma 3, the random variables u and v are not necessarily zero-mean. Consequently, the expectations $\mathbb{E}\{\operatorname{sgn}\{u\}\operatorname{sgn}\{v\}\}\$ and $\mathbb{E}\{u\operatorname{sgn}\{v\}\}\$ cannot be computed via Price's theorem as in [11], [13], [14]. We can compute Q_4 by making the following identifications:

$$u \leftarrow [\tilde{\boldsymbol{w}}_n]_i$$
 (30)

$$v \leftarrow [\boldsymbol{w}^{\star} + \tilde{\boldsymbol{w}}_n]_j \tag{31}$$

with

$$u_u \leftarrow \mathbb{E}\{[\tilde{\boldsymbol{w}}_n]_i\}$$
 (32)

$$\mu_v \leftarrow \mathbb{E}\{[\boldsymbol{w}^\star + \tilde{\boldsymbol{w}}_n]_j\}$$
(33)

The variances and covariance σ_u^2 , σ_v^2 and ρ_{uv} have the same expressions as in (25)–(27) because they are invariant to a constant shift.

The term Q_5 is then given by:

$$\boldsymbol{Q}_5 = \boldsymbol{R}_x \, \boldsymbol{Q}_4. \tag{34}$$

Replacing Q_1 to Q_5 into (12) we can characterize the secondorder behavior of the ZA-LMS algorithm.

IV. EXPERIMENT VALIDATION

In this section, we present simulation examples to validate our models. Consider an unknown system of order L = 17with weights defined by:

$$\boldsymbol{w}^{\star} = [0.8, \ 0.5, \ 0.3, \ 0.1, \ 0.05, \ \boldsymbol{0}_{5}^{\top}, \\ -0.05, \ -0.1, \ -0.3, \ -0.5, \ -0.8]^{\top},$$
(35)

The input signal was a first-order AR process defined as follows: $x_n = 0.6 x_{n-1} + \nu_n$, with ν_n an i.i.d. zero-mean Gaussian variable with variance $\sigma_{\nu}^2 = 0.64$ (so that $\sigma_x^2 = 1$). The additive noise z_n was zero-mean i.i.d. Gaussian with variance $\sigma_z^2 = 0.01$. The adaptive weights were initialized to $w_0 = \mathbf{0}_L$. The step size was set to $\mu = 0.01$, and the regularization parameter was set to $\lambda = 0.01$ and 0.001. Note that our objective was not to evaluate the performance of the algorithm but to evaluate the accuracy of our models. The tuning parameters were thus set arbitrarily.

Before testing our models, we show histograms in Fig. 2 to support assumption A2 used in the analysis. They were built from 5000 samples of bivariate vectors $[[\tilde{\boldsymbol{w}}_n]_3, [\tilde{\boldsymbol{w}}_n]_8]$ at n = 100 and n = 800, and $[[\tilde{\boldsymbol{w}}_n]_9, [\tilde{\boldsymbol{w}}_n]_{10}]$ at n = 800, for $\mu = 0.01$ and $\lambda = 0.001$. It can be observed that they have Gaussian-like profiles, in particular at steady-state n = 800.

The accuracy of our models are illustrated in Fig. 1. Simulation results were obtained by averaging over 500 runs. Fig. 1(a) illustrates the mean weight behavior (6) of ZA-LMS with $\lambda = 0.01$. The simulated curves (blue) and theoretical curves (6)-(7) (red) are superimposed. Fig. 1(b) shows the MSE and EMSE learning curves (11), along with the theoretical curves (11)-(17) obtained in this paper and in [8] for comparison purpose. For transient MSE, both models are generally consistent with the simulated curve, while a zoom-in on the interval [900, 1000] shows that our model has a better consistency. For transient EMSE, our model is consistent with the simulated results while the model proposed in [8] has a large bias. We thus can conclude that in the MSE curve, the inaccuracy of the model in [8] is hidden by the offset σ_z^2 in (11) that is much higher than EMSE. Fig. 1(c) shows a similar result with $\lambda = 0.001$.

V. CONCLUSION

In this paper, we derived analytical models to characterize the transient behavior of ZA-LMS in the mean and meansquare sense. Simulations illustrated the consistency between



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Fig. 1. Model validation and comparison with model obtained in [8].

Monte Carlo learning curves and the theoretical findings, as well as the improved accuracy of our model compared with previous works. Extension to more general update relation such as the reweighted ZA-LMS algorithm will be considered in future works.

APPENDIX A: PROOF OF LEMMA 1

Consider a Gaussian random variable $u \sim \mathcal{N}(\mu, \sigma^2)$. We have:

$$\mathbb{E}\{\operatorname{sgn}\{u\}\} = \int_{-\infty}^{0} -\mathcal{N}(\mu, \sigma^2) du + \int_{0}^{+\infty} \mathcal{N}(\mu, \sigma^2) du \qquad (36)$$
$$= 1 - 2\phi(-\mu/\sigma).$$

APPENDIX B: PROOF OF LEMMA 2

Consider jointly Gaussian variables u and v. We have:

$$\mathbb{E}\{\operatorname{sgn}\{u\}\operatorname{sgn}\{v\}\} = \int_{-\infty}^{0} \int_{-\infty}^{0} \mathcal{N}([u, v]^{\top}, \boldsymbol{\Sigma}_{uv}) \, du \, dv + \int_{0}^{+\infty} \int_{0}^{+\infty} \mathcal{N}([u, v]^{\top}, \boldsymbol{\Sigma}_{uv}) \, du \, dv - \int_{-\infty}^{0} \int_{0}^{+\infty} \mathcal{N}([u, v]^{\top}, \boldsymbol{\Sigma}_{uv}) \, du \, dv - \int_{0}^{+\infty} \int_{-\infty}^{0} \mathcal{N}([u, v]^{\top}, \boldsymbol{\Sigma}_{uv}) \, du \, dv.$$
(37)

With simple algebraic manipulations and using the definition of CDF of the multivariate Gaussian distribution, we write:

$$\mathbb{E}\{\operatorname{sgn}\{u\}\operatorname{sgn}\{v\}\}$$

= $\Phi(\mathbf{0}_2, [\mu_u, \mu_v]^\top, \boldsymbol{\Sigma}_{uv}) + \Phi(\mathbf{0}_2, -[\mu_u, \mu_v]^\top, \boldsymbol{\Sigma}_{uv})$ (38)
- $\Phi(\mathbf{0}_2, [\mu_u, -\mu_v]^\top, \overline{\boldsymbol{\Sigma}}_{uv}) - \Phi(\mathbf{0}_2, [-\mu_u, \mu_v]^\top, \overline{\boldsymbol{\Sigma}}_{uv})$

where
$$\overline{\Sigma}_{uv} = \Sigma_{uv} \circ \begin{bmatrix} 1 & -1 \\ -1 & 1 \end{bmatrix}$$
.

APPENDIX C: PROOF OF LEMMA 3

Consider jointly Gaussian variables u and v. We have:

$$\mathbb{E}\{u\operatorname{sgn}\{v\}\} = \int_{-\infty}^{+\infty} \operatorname{sgn}\{v\} \left(\int_{-\infty}^{+\infty} u \,\mathcal{N}([u,v]^{\top}, \boldsymbol{\Sigma}_{uv}) \,du\right) dv.$$
⁽³⁹⁾

For ease of presentation, we write:

$$\boldsymbol{\Sigma}_{uv}^{-1} = \begin{bmatrix} a & c \\ c & b \end{bmatrix}.$$
 (40)

Since Σ_{uv} is positive definite, we have a > 0 and $ab - c^2 > 0$. This remark is important since $b - \frac{c^2}{a}$ will denote below the variance of a Gaussian distribution. The integral (39) can be expressed via the above quantities:

$$(39) = \frac{1}{2\pi\sqrt{|\boldsymbol{\Sigma}_{uv}|}} \int_{-\infty}^{+\infty} \left\{ \operatorname{sgn}\{v\} \exp\left[-\frac{1}{2}(b - \frac{c^2}{a})(v - \mu_v)^2\right] \times \left(\sqrt{\frac{2\pi}{a}} \int_{-\infty}^{+\infty} u \,\mathcal{N}(\mu_u - \frac{c}{a}(v - \mu_v), \frac{1}{a}) \,du\right) \right\} dv$$

$$= \frac{1}{\sqrt{2\pi a |\boldsymbol{\Sigma}_{uv}|}} \int_{-\infty}^{+\infty} \left\{ \operatorname{sgn}\{v\} \exp\left[-\frac{1}{2}(b - \frac{c^2}{a})(v - \mu_v)^2\right] \times \left(\mu_u - \frac{c}{a}(v - \mu_v)\right) \right\} dv$$

$$= \frac{1}{\sqrt{2\pi a |\boldsymbol{\Sigma}_{uv}|}} \times \left\{ \int_{-\infty}^{+\infty} \left(\mu_u + \frac{c}{a}\mu_v\right) \operatorname{sgn}\{v\} \exp\left[-\frac{1}{2}\delta(v - \mu_v)^2\right] dv$$

$$- \int_{-\infty}^{+\infty} \frac{c}{a} v \operatorname{sgn}\{v\} \exp\left[-\frac{1}{2}\delta(v - \mu_v)^2\right] dv \right\}$$

$$(41)$$

where $\delta = b - \frac{c^2}{a} > 0$. The second equality comes from the fact that the second integral in the first equality is the mean value of $\mathcal{N}(\mu_u - \frac{c}{a}(v - \mu_v), \frac{1}{a})$. Using Lemma 1, we get:

$$\int_{-\infty}^{+\infty} \left(\mu_u + \frac{c}{a}\mu_v\right) \operatorname{sgn}\{v\} \exp\left[-\frac{1}{2}\delta(v-\mu_v)^2\right] dv$$

$$= \left(\mu_u + \frac{c}{a}\mu_v\right) \sqrt{\frac{2\pi}{\delta}} \left[1 - 2\phi\left(\mu_v\sqrt{\delta}\right)\right].$$
(42)

Finally, by considering the mean value of the folded Gaussian distribution [15], we have:

$$\begin{aligned} \int_{-\infty}^{+\infty} \frac{c}{a} v \operatorname{sgn}\{v\} \exp\left[-\frac{1}{2}\delta(v-\mu_{v})^{2}\right] dv \\ &= \frac{c}{a} \sqrt{\frac{2\pi}{\delta}} \int_{-\infty}^{+\infty} |v| \frac{1}{\sqrt{2\pi}\sqrt{1/\delta}} \exp\left[-\frac{(v-\mu_{v})^{2}}{2(1/\delta)}\right] dv \\ &= \frac{c}{a} \sqrt{\frac{2\pi}{\delta}} \left[\sqrt{\frac{2}{\delta\pi}} \exp\left(-\frac{1}{2}\mu_{v}^{2}\delta\right) + \mu_{v} \left(1 - 2\phi\left(-\mu_{v}\sqrt{\delta}\right)\right)\right] \end{aligned}$$
(43)

Substituting the results in (42) and (43) into (41) yields the value of $\mathbb{E}\{u \operatorname{sgn}\{v\}\}$.

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