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Stochastic Behavior Analysis of the Gaussian Kernel Least-Mean-Square Algorithm

Wemerson D. Parreira, José Carlos M. Bermudez, Cédric Richard and Jean-Yves Tourneret

Abstract—The kernel least-mean-square (KLMS) algorithm is a popular algorithm in nonlinear adaptive filtering due to its simplicity and robustness. In kernel adaptive filters, the statistics of the input to the linear filter depends on the parameters of the kernel employed. Moreover, practical implementations require a finite nonlinearity model order. A Gaussian KLMS has two design parameters, the step size and the Gaussian kernel bandwidth. Thus, its design requires analytical models for the algorithm behavior as a function of these two parameters. This paper studies the steady-state behavior and the transient behavior of the Gaussian KLMS algorithm for Gaussian inputs and a finite order nonlinearity model. In particular, we derive recursive expressions for the mean-weight-error vector and the mean-square-error. The model predictions show excellent agreement with Monte Carlo simulations. Nonlinear systems can be characterized by representations ranging from higher-order statistics to series expansion methods [1]. Nonlinear system identification methods based on reproducing kernel Hilbert spaces (RKHS) have gained popularity over the last decades [2], [3]. According to this rule, the function \( \psi(\cdot) \) that minimizes the squared estimation error \( \sum_{n=1}^{N} (d(n) - \psi(u(n)))^2 \), given \( N \) input vectors \( u(n) \) and desired outputs \( d(n) \), can be written as the kernel expansion \( \psi(\cdot) = \sum_{n=1}^{N} \alpha_n \kappa(\cdot, u(n)) \). This reduces the problem to determining \( \alpha = [\alpha_1, \ldots, \alpha_N]^T \) that minimizes \( \|d - K\alpha\|^2 \), where \( K \) is the Gram matrix with \( (n, \ell) \)th entry \( \kappa(u(n), u(\ell)) \), and \( d = [d(1), \ldots, d(N)]^T \). Since the order of the model is equal to the number \( N \) of available data \( u(n) \), this approach cannot be considered for online applications. To overcome this barrier, one can focus on finite-order models

\[
\psi(\cdot) = \sum_{j=1}^{M} \alpha_j \kappa(\cdot, u(\omega_j))
\]

where \( \omega_j, j = 1, \ldots, M \) form a subset of \( \{1, 2, \ldots, N\} \) corresponding to the time indexes of the \( M < N \) input vectors chosen to build the \( M \)th-order model (1). The \( M \) kernel functions \( \kappa(\cdot, u(\omega_j)) \) form the dictionary. In [4], the authors present an overview of the existing techniques to select the kernel functions in (1), an example of which is the approximate linear dependence (ALD) criterion [5]. It consists of including a kernel function \( \kappa(\cdot, u(\ell)) \) in the dictionary if it satisfies

\[
\min_{\gamma} \|\kappa(\cdot, u(\ell)) - \sum_j \gamma_j \kappa(\cdot, u(\omega_j))\|_2^2 > \epsilon_0
\]

where \( \epsilon_0 \) is a parameter determining model sparsity level. To control the model order with reduced computational complexity, the coherence-based sparsification rule has also been considered [6], [7]. According to this rule, the kernel \( \kappa(\cdot, u(\ell)) \) is inserted into the dictionary if

\[
\max_j |\kappa(\cdot, u(\ell))| < \epsilon_0
\]

with \( \epsilon_0 \) a parameter determining the dictionary coherence. It was shown in [7] that the dictionary dimension determined under rule (3) is finite. For the rest of the paper, we shall assume that the dictionary size \( M \) is known, fixed and finite.

It is well known that a nonlinear adaptive filtering problem with input signal in \( \mathcal{U} \) can be solved using a linear adaptive filter [4]. The linear adaptive filter input is a nonlinear mapping of \( \mathcal{U} \) to an Hilbert space \( \mathcal{H} \) possessing a reproducing kernel. The theory outlined above shows that the order of the linear adaptive filter can be finite if a proper input sparsification rule is employed, even if the dimensionality of the transformed input in

I. INTRODUCTION

M ANY practical applications (e.g., in communications and bioengineering) require nonlinear signal processing. Nonlinear systems can be characterized by representations ranging from higher-order statistics to series expansion methods [1]. Nonlinear system identification methods based on reproducing kernel Hilbert spaces (RKHS) have gained popularity over the last decades [2], [3]. More recently, kernel adaptive filtering has been recognized as an appealing solution to the nonlinear adaptive filtering problem, as working in RKHS allows the use of linear structures to solve nonlinear estimation problems. See [4] for an overview. The block diagram of a kernel-based adaptive system identification problem is shown in Fig. 1. Here, \( \mathcal{U} \) is a compact subspace of \( \mathbb{H}^2 \). \( \kappa: \mathcal{U} \times \mathcal{U} \rightarrow \mathbb{H} \) is a reproducing kernel, \( (\mathcal{H}, \langle \cdot, \cdot \rangle) \) is the induced RKHS with its inner product and \( z(n) \) is a zero-mean additive noise uncorrelated with any other signal. The representer theorem [2] states that the function \( \psi(\cdot) \) that minimizes the squared estimation error \( \sum_{n=1}^{N} (d(n) - \psi(u(n)))^2 \), given \( N \) input vectors \( u(n) \) and desired outputs \( d(n) \), can be written as the kernel expansion \( \psi(\cdot) = \sum_{n=1}^{N} \alpha_n \kappa(\cdot, u(n)) \). This reduces the problem to determining \( \alpha = [\alpha_1, \ldots, \alpha_N]^T \) that minimizes \( \|d - K\alpha\|^2 \), where \( K \) is the Gram matrix with \( (n, \ell) \)th entry \( \kappa(u(n), u(\ell)) \), and \( d = [d(1), \ldots, d(N)]^T \). Since the order of the model is equal to the number \( N \) of available data \( u(n) \), this approach cannot be considered for online applications. To overcome this barrier, one can focus on finite-order models

\[
\psi(\cdot) = \sum_{j=1}^{M} \alpha_j \kappa(\cdot, u(\omega_j))
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where \( \omega_j, j = 1, \ldots, M \) form a subset of \( \{1, 2, \ldots, N\} \) corresponding to the time indexes of the \( M < N \) input vectors chosen to build the \( M \)th-order model (1). The \( M \) kernel functions \( \kappa(\cdot, u(\omega_j)) \) form the dictionary. In [4], the authors present an overview of the existing techniques to select the kernel functions in (1), an example of which is the approximate linear dependence (ALD) criterion [5]. It consists of including a kernel function \( \kappa(\cdot, u(\ell)) \) in the dictionary if it satisfies

\[
\min_{\gamma} \|\kappa(\cdot, u(\ell)) - \sum_j \gamma_j \kappa(\cdot, u(\omega_j))\|_2^2 > \epsilon_0
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where \( \epsilon_0 \) is a parameter determining model sparsity level. To control the model order with reduced computational complexity, the coherence-based sparsification rule has also been considered [6], [7]. According to this rule, the kernel \( \kappa(\cdot, u(\ell)) \) is inserted into the dictionary if

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Fig. 1. Kernel-based adaptive system identification.

\( \mathcal{H} \) is infinite as in the case of Gaussian kernel. Algorithms developed using these ideas include the kernel least-mean-square (KLMS) algorithm [8], [9], the kernel recursive-least-square (KRLS) algorithm [5], the kernel-based normalized least-mean-square (KNLMS) algorithm and the affine projection (KAPA) algorithm [6], [10], [7]. See also the monograph [4]. In addition to the choice of the usual linear adaptive filter parameters, designing kernel adaptive filters requires the choice of the kernel and its parameters. Moreover, using the finite order model in (1) implies that the adaptive algorithm behavior cannot be studied as the behavior of the algorithm presented in [4, (2.17)], which is a regular LMS algorithm operating in the RKHS. Choosing the algorithm and nonlinearity model parameters to achieve a prescribed performance is still an open issue, and requires an extensive analysis of the algorithm stochastic behavior. Our work brings a new contribution to the discussion about kernel-based adaptive filtering by providing the first convergence analysis of the KLMS algorithm with Gaussian kernel.

The paper is structured as follows. In Section II, we derive recursive expressions for the mean weight error vector and the mean-square error (MSE) for Gaussian inputs. In Section III, we define analytical models for the transient behavior of the first and second-order moments of the adaptive weights. Section IV studies the algorithm convergence properties. Stability conditions and a steady-state behavior model are derived which allow the algorithm design for prescribed convergence speed and quality of estimate. In Section V, we use the analysis results to establish design guidelines. Section VI presents design examples which validate the theoretical analysis and illustrate its application. The model predictions show excellent agreement with Monte Carlo simulations both in transient and steady state.

II. MEAN SQUARE ERROR ANALYSIS

Consider the nonlinear system identification problem shown in Fig. 1, and the finite-order model (1) based on the Gaussian kernel

\[
\kappa(\mathbf{u}, \mathbf{u}') = \exp \left( -\frac{|\mathbf{u} - \mathbf{u}'|^2}{2\xi^2} \right) \tag{4}
\]

where \( \xi \) is the kernel bandwidth. The environment is assumed stationary, meaning that \( \psi(\mathbf{u}(n)) \) is stationary for \( \mathbf{u}(n) \) stationary. This assumption is satisfied by several nonlinear systems used to model practical situations, such as memoryless, Wiener and Hammerstein systems. System inputs are zero-mean, independent, and identically distributed Gaussian \((q \times 1)\) vectors \( \mathbf{u}(n) \) so that \( E[\mathbf{u}(n - i)\mathbf{u}^\top(n - j)] = 0 \) for \( i \neq j \). The components of the input vector \( \mathbf{u}(n) \) can, however, be correlated. Let \( \mathbf{R}_{\mathbf{uu}} = E[\mathbf{u}(n)\mathbf{u}^\top(n)] \) denote their autocorrelation matrix.

For a dictionary of size \( M \), let \( \kappa_\omega(n) \) be the vector of kernels at time \( n > M \), that is

\[
\kappa_\omega(n) = [\kappa(\mathbf{u}(n), \mathbf{u}(\omega_1)), \ldots, \kappa(\mathbf{u}(n), \mathbf{u}(\omega_M))]^\top \tag{5}
\]

where \( \mathbf{u}(\omega_i) \) is the \( i \)th element of the dictionary, with \( \mathbf{u}(\omega_i) \neq \mathbf{u}(n) \) for \( i = 1, \ldots, M \). Here we consider that the vectors \( \mathbf{u}(\omega_i) \), \( i = 1, \ldots, M \) may change at each iteration following some dictionary updating schedule. The only limitation imposed in the following analysis is that \( \kappa_\omega(n) \neq \mathbf{u}(\omega_i) \), \( i \neq j \) so that the dictionary vectors which are arguments of different entries of \( \kappa_\omega(n) \) are statistically independent. To keep the notation simple, however, we will not show explicitly the dependence of \( \omega_i \) on \( n \) and represent \( \kappa(\mathbf{u}(n), \mathbf{u}(\omega_i)) \) as \( \kappa(\omega_i) \) for all \( i \).

From Fig. 1 and model (1), the estimated system output is

\[
\hat{d}(n) = \mathbf{a}(n)^\top \kappa_\omega(n) \tag{6}
\]

with \( \mathbf{a}(n) = [\alpha_1(n), \ldots, \alpha_M(n)]^\top \). The corresponding estimation error is defined as

\[
e(n) = \hat{d}(n) - d(n). \tag{7}
\]

Squaring both sides of (7) and taking the expected value leads to the MSE

\[
J_{\text{max}}(n) = E[e^2(n)] = E[d^2(n)] - 2p_{\text{ed}} \mathbf{a}(n) + \mathbf{a}(n)^\top \mathbf{R}_{\text{xx}} \mathbf{a}(n) \tag{8}
\]

where \( \mathbf{R}_{\text{xx}} = E[\kappa_\omega(n)\kappa_\omega^\top(n)] \) is the correlation matrix of the kernelized input, and \( p_{\text{ed}} = E[d(n)\kappa_\omega(n)] \) is the cross-correlation vector between \( \kappa_\omega(n) \) and \( d(n) \). It is shown in Appendix A that \( \mathbf{R}_{\text{xx}} \) is positive definite. Thus, the optimum weight vector is given by

\[
\mathbf{a}_{\text{opt}} = \mathbf{R}_{\text{xx}}^{-1} p_{\text{ed}} \tag{9}
\]

and the corresponding minimum MSE is

\[
J_{\text{min}} = E[d^2(n)] - p_{\text{ed}}^\top \mathbf{R}_{\text{xx}}^{-1} p_{\text{ed}}. \tag{10}
\]

These are the well-known expressions of the Wiener solution and minimum MSE, where the input signal vector has been replaced by the kernelized input vector. Determining the optimum \( \mathbf{a}_{\text{opt}} \) requires the determination of the covariance matrix

1If the dictionary size \( M \) is adapted online, assume that \( n \) is sufficiently large so that the size \( M \) does not increase anymore.
$R_{\kappa\kappa}$, given the statistical properties of $u(n)$ and the reproducing kernel.

Before closing this section, let us evaluate the correlation matrix $R_{\kappa\kappa}$. Its entries are given by

$$
\begin{align*}
| R_{\kappa\kappa} |_{ij} &= \begin{cases} 
E\{\kappa^2(u(n), u(\omega_i))\}, & i = j \\
E\{\kappa(u(n), u(\omega_i))\} r(\kappa(u(n), u(\omega_j))), & i \neq j
\end{cases}
\end{align*}
$$

with $1 \leq i, j \leq M$. Note that $R_{\kappa\kappa}$ remains time-invariant even if the dictionary is updated at each iteration, as $u(n)$ is stationary and $u(\omega_i)$ and $u(\omega_j)$ are statistically independent for $i \neq j$.

Let us introduce the following notations

$$
\| \cdot \| \text{ is the } L_2 \text{ norm and}
$$

$$
\begin{align*}
y_2 &= \left( u^T(n) u(\omega_i) \right)^T \\
y_3 &= \left( u^T(n) u^T(\omega_i) u^T(\omega_j) \right)^T
\end{align*}
$$

and

$$
Q_2 = \begin{pmatrix} I & -I \\
-I & I \end{pmatrix} \quad Q_3 = \begin{pmatrix} 2I & -I & -I \\
-I & I & 0 \\
-I & 0 & I \end{pmatrix}
$$

where $I$ is the $(q \times q)$ identity matrix and $O$ is the $(q \times q)$ null matrix. From [11, p. 100], we know that the moment generating function of a quadratic form $z = y^T Q y$, where $y$ is a zero-mean Gaussian vector with covariance matrix $R_y$, is given by

$$
\psi_z(s) = E[e^{s^T z}] = \det\{I - 2sQ R_y\}^{-1/2}.
$$

Making $s = -1/(2\xi^2)$ in (15), we find that the $(i,j)$th element of $R_{\kappa\kappa}$ is given by

$$
| R_{\kappa\kappa} |_{ij} = \begin{cases} 
\text{det}\{I_3 + 2Q_3 R_z / \xi^2\}^{-1/2}, & i = j \\
\text{det}\{I_3 + Q_3 R_z / \xi^2\}^{-1/2}, & i \neq j
\end{cases}
$$

with $1 \leq i, j \leq M$. The main diagonal entries $| R_{\kappa\kappa} |_{ii}$ are all equal to $r_{\text{nd}}$ and the off-diagonal entries $| R_{\kappa\kappa} |_{ij}$ are all equal to $r_{\text{nd}}$ since $u(\omega_i)$ and $u(\omega_j)$ are i.i.d. In (16), $R_z$ is the $(\ell_2 \times \ell_2)$ correlation matrix of vector $y$, $I_i$ is the $(\ell_2 \times \ell_2)$ identity matrix, and $\text{det}\{\cdot\}$ denotes the determinant of a matrix. Finally, note that matrix $R_z$ is block-diagonal with $R_{\kappa\kappa}$ along its diagonal.

### III. GAUSSIAN KLMS ALGORITHM: TRANSIENT BEHAVIOR ANALYSIS

The KLMS weight-update equation for the system presented in Fig. 1 is [4]

$$
\alpha(n+1) = \alpha(n) + \eta e(n) \omega(n).
$$

Defining the weight-error vector $w(n) = \alpha(n) - \alpha_{\text{opt}}$ leads to the weight-error vector update equation

$$
\begin{align*}
w(n+1) &= w(n) + \eta e(n) \omega(n). \\
\text{From (6) and (7), and the definition of } w(n), \text{ the error equation is given by}
\end{align*}
$$

$$
e(n) = d(n) - \omega(n) - \omega(n)^T \alpha_{\text{opt}}
$$

and the optimal estimation error is

$$
e_0(n) = d(n) - \omega(n)^T \alpha_{\text{opt}}.
$$

Substituting (19) into (18) yields

$$
\begin{align*}
u(n+1) &= v(n) + \eta \omega(n) \omega(n) - \eta \omega(n)^T \omega(n) \omega(n) - \eta \omega(n)^T \omega(n) \alpha_{\text{opt}}.
\end{align*}
$$

### A. Simplifying Statistical Assumptions

Simplifying assumptions are required in order to make the study of the stochastic behavior of $u(n)$ mathematically feasible. The statistical assumptions required in different parts of the analysis are the following:

- **A1:** $\omega(n)$ is statistically independent of $u(n)$.
  This assumption is justified in detail in [12] and has been successfully employed in several adaptive filter analyses. It is called here for further reference “modified independence assumption” (MIA). This assumption has been shown in [12] to be less restrictive than the classical independence assumption [13].

- **A2:** The finite-order model provides a close enough approximation to the infinite-order model with minimum MSE, so that $E[e_0(n)^2] \approx 0$.

- **A3:** $e_0(n)$ and $\omega(n)$ are uncorrelated. This assumption is also supported by the arguments supporting the MIA (A1) [12].

### B. Mean Weight Behavior

Taking the expected value of both sides of (21) and using the MIA (A1) yields

$$
E\{u(n+1)\} = (I - \eta R_{\kappa\kappa}) E\{u(n)\}
$$

which is the LMS mean weight behavior for an input vector $\omega(n)$.

### C. Mean-Square Error

Using (19) and the MIA (A1), the second-order moments of the weights are related to the MSE through [13]

$$
J_{\text{mse}}(n) = J_{\text{min}} + \text{tr}\{R_{\kappa\kappa} C_{\text{e}}(n)\}
$$

where $C_{\text{e}}(n) = E\{w(n) w^H(n)\}$ is the autocorrelation matrix of $w(n)$ and $J_{\text{min}} = E\{e_0^2(n)\}$ the minimum MSE. The study of the MSE behavior (23) requires a model for $C_{\text{e}}(n)$. This model is highly affected by the transformation imposed on the input signal $u(n)$ by the kernel. An analytical model for the behavior of $C_{\text{e}}(n)$ is derived in the next subsection.

### D. Second-Order Moment Behavior

Using (20) and (21), the weight-error vector update becomes

$$
\begin{align*}
u(n+1) &= v(n) + \eta e_0(n) \omega(n) \\
-\eta \omega(n)^T \omega(n) \omega(n) + \eta \omega(n)^T \omega(n) \alpha_{\text{opt}}.
\end{align*}
$$

From (6) and (7), and the definition of $w(n)$, the error equation is given by

$$
e(n) = d(n) - \omega(n) - \omega(n)^T \alpha_{\text{opt}}
$$

and the optimal estimation error is

$$
e_0(n) = d(n) - \omega(n)^T \alpha_{\text{opt}}.
$$

Substituting (19) into (18) yields

$$
\begin{align*}
u(n+1) &= v(n) + \eta d(n) \omega(n) - \eta \omega(n)^T \omega(n) \omega(n) - \eta \omega(n)^T \omega(n) \alpha_{\text{opt}}.
\end{align*}
$$

From (6) and (7), and the definition of $w(n)$, the error equation is given by

$$
e(n) = d(n) - \omega(n) - \omega(n)^T \alpha_{\text{opt}}
$$

and the optimal estimation error is

$$
e_0(n) = d(n) - \omega(n)^T \alpha_{\text{opt}}.
$$

Substituting (19) into (18) yields

$$
\begin{align*}
v(n+1) &= v(n) + \eta d(n) \omega(n) - \eta \omega(n)^T \omega(n) \omega(n) - \eta \omega(n)^T \omega(n) \alpha_{\text{opt}}.
\end{align*}
$$

From (6) and (7), and the definition of $w(n)$, the error equation is given by

$$
e(n) = d(n) - \omega(n) - \omega(n)^T \alpha_{\text{opt}}
$$

and the optimal estimation error is

$$
e_0(n) = d(n) - \omega(n)^T \alpha_{\text{opt}}.
$$

Substituting (19) into (18) yields

$$
\begin{align*}
v(n+1) &= v(n) + \eta \omega(n) - \eta \omega(n)^T \omega(n) \omega(n) - \eta \omega(n)^T \omega(n) \alpha_{\text{opt}}.
\end{align*}
$$
Post-multiplying (24) by its transpose, and taking the expected value, leads to

\[
C_v(n+1) = C_v(n) - \eta E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{v}(n) \mathbf{v}^\top(n)\}
\]

\[
- \eta E\{\mathbf{v}(n) \mathbf{v}^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\} + \eta^2 E\{\mathbf{v}_0(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

\[
+ \eta E\{v_0(n) \mathbf{k}_\omega(n) \mathbf{v}^\top(n)\}
\]

\[
- \eta^2 E\{v_0(n) \mathbf{k}_\omega(n) \mathbf{v}^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

Using the MIA (A1), the first two expected values are given by

\[
E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{v}(n) \mathbf{v}^\top(n)\} = \mathbf{R}_{\mathbf{k}_\omega} C_v(n)
\]

\[
E\{\mathbf{v}(n) \mathbf{v}^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\} = C_v(n) \mathbf{R}_{\mathbf{k}_\omega}.
\]

Using assumptions A2 and A3 the third expected value is given by

\[
E\{E_0(n) \mathbf{k}_\omega(n) \mathbf{v}^\top(n)\} \approx E\{E_0(n) \mathbf{k}_\omega(n)\} E\{\mathbf{v}^\top(n)\}
\]

\[
= \mathbf{R}_{\mathbf{k}_\omega} J_{\mathbf{v},\mathbf{v}}.
\]

The fourth and the sixth expected values can be approximated using the MIA (A1), that is,

\[
E\{E_0(n) \mathbf{k}_\omega(n) \mathbf{v}^\top(n)\} \approx E\{E_0(n) \mathbf{k}_\omega(n)\} E\{\mathbf{v}^\top(n)\}
\]

\[
= 0
\]

since \(E\{E_0(n) \mathbf{k}_\omega(n)\} = 0\) by the orthogonality principle [13].

Evaluation of the fifth and seventh expected values requires further simplifications for mathematical tractability. A reasonable approximation that preserves the effect of \(\mathbf{k}_\omega(n)\) up to its second-order moments is to assume that \(E_0(n) \mathbf{k}_\omega(n)\) and \(\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\) are uncorrelated. Under both this approximation and MIA (A1),

\[
E\{E_0(n) \mathbf{k}_\omega(n) \mathbf{v}^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

\[
= E\{\mathbf{v}^\top(n)\} E\{E_0(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\} = 0
\]

where the equality to zero is due to the orthogonality principle.

Using (26)–(29) in (25) yields

\[
C_v(n+1) \approx C_v(n) - \eta \mathbf{R}_{\mathbf{k}_\omega} C_v(n) + \mathbf{C}_v(n) \mathbf{R}_{\mathbf{k}_\omega} + \eta^2 \mathbf{T}(n) + \eta^2 \mathbf{R}_{\mathbf{k}_\omega} J_{\mathbf{v},\mathbf{v}}
\]

(30a)

with

\[
\mathbf{T}(n) = E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{v}(n) \mathbf{v}^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}.
\]

(30b)

Evaluation of expectation (30b) is an important step in the analysis. In the classical LMS analysis [14], the input signal is assumed zero-mean Gaussian. Then the expectation in (30b) can be approximated using the moment factoring theorem for Gaussian variates. In the present analysis, as \(\mathbf{k}_\omega(n)\) is a non-linear transformation of a quadratic function of the Gaussian input vector \(\mathbf{u}(n)\), it is neither zero-mean nor Gaussian.

Using the MIA (A1) to determine the \((i, j)\)th element of \(\mathbf{T}(n)\) in (30b) yields

\[
[T(n)]_{ij} \approx \sum_{\ell=1}^{M} \sum_{p=1}^{M} E\{\mathbf{v}(n) \mathbf{v}^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\} \times [\mathbf{v}(n)]_{\ell p}
\]

where \(\mathbf{k}_\omega(n) = \mathbf{\sigma}(\mathbf{u}(n), \mathbf{w}(n))\). Each of the moments in (31) can be determined using (15). Depending on \(i, j, \ell, p\), we have five different moments \(\mu_k\), \(k = 1, \ldots, 5\) to evaluate.

\[
\mu_1 := E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

\[
\mu_2 := E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

\[
\mu_3 := E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

\[
\mu_4 := E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

\[
\mu_5 := E\{\mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n) \mathbf{k}_\omega(n) \mathbf{k}_\omega^\top(n)\}
\]

(31)

Denoting \(y_\mathbf{v} = (\mathbf{u}(n) \mathbf{u}^\top(n))\), \(y_\mathbf{k} = (\mathbf{w}(n) \mathbf{w}^\top(n))\),

\[
\mu_1 = |\mathbf{I} + Q_2 \mathbf{R}_{\mathbf{v}}|^{-1/2}
\]

(32)

\[
\mu_2 = |\mathbf{I} + Q_2 \mathbf{R}_{\mathbf{v}}|^{-1/2}
\]

(33)

where

\[
Q_2 = \begin{pmatrix}
4I & -3I & -I \\
-3I & 3I & O \\
-I & O & I
\end{pmatrix}
\]

(34)

Denoting \(y_3 = (\mathbf{w}(n) \mathbf{u}^\top(n))\),

\[
\mu_3 = |\mathbf{I} + Q_3 \mathbf{R}_{\mathbf{v}}|^{-1/2}
\]

(35)

Denoting \(y_4 = (\mathbf{w}(n) \mathbf{u}^\top(n))\),

\[
\mu_4 = |\mathbf{I} + Q_4 \mathbf{R}_{\mathbf{v}}|^{-1/2}
\]

(36)

where

\[
Q_3 = \begin{pmatrix}
4I & 2I & -I & -I \\
-2I & 2I & O & O \\
-I & O & I & O \\
-I & O & O & I
\end{pmatrix}
\]

(37)

Denoting \(y_5 = (\mathbf{w}(n) \mathbf{w}^\top(n))\),

\[
\mu_5 = |\mathbf{I} + Q_5 \mathbf{R}_{\mathbf{v}}|^{-1/2}
\]

(38)
where
\[
Q_5 = \begin{pmatrix}
-4I & -I & -I & -I & -I \\
-I & I & 0 & 0 & 0 \\
-I & 0 & I & 0 & 0 \\
-I & 0 & 0 & I & 0 \\
-I & 0 & 0 & 0 & I \\
\end{pmatrix}
\] (39)

Using these moments, the elements of \(T(n)\) are finally given by
\[
[T(n)]_{ii} = \mu_3 \mathbf{C}_v(n)_{ii} + \mu_4 \sum_{p \neq 1} \mathbf{C}_v(n)_{ip} + \mu_5 \sum_{p \neq 1} \mathbf{C}_v(n)_{ip} + \mu_3 \mathbf{C}_v(n)_{ii} \tag{40}
\]
and, for \(j \neq i\),
\[
[T(n)]_{ij} = \sum_{\ell \neq 1} \left\{ 2 \mu_2 \mathbf{C}_v(n)_{\ell j} + \mu_3 \mathbf{C}_v(n)_{\ell j} \right\} + \mu_4 \mathbf{C}_v(n)_{i j} + \mu_5 \sum_{p \neq 1} \mathbf{C}_v(n)_{ip} \tag{41}
\]
which completes the evaluation of \(T(n)\) in (30b). Substituting this result into (30a) yields the following recursive expressions for the entries of the autocorrelation matrix \(\mathbf{C}_v(n)\):
\[
\mathbf{C}_v(n + 1)_{ii} = (1 - 2\eta r_{mod} + \eta^2 \mu_1) \mathbf{C}_v(n)_{ii} + \eta^2 \mu_3 \sum_{\ell = 1}^{M} \mathbf{C}_v(n)_{i\ell} + (2\eta^2 \mu_2 - 2\eta r_{mod}) \sum_{\ell = 1}^{M} \mathbf{C}_v(n)_{i\ell} + \eta^2 \mu_4 \sum_{\ell = 1}^{M} \mathbf{C}_v(n)_{i\ell} + \eta^2 \mu_5 \sum_{p = 1}^{M} \mathbf{C}_v(n)_{ip} + \eta^2 r_{mod} J_{\text{mix}} \tag{42}
\]
and, for \(j \neq i\),
\[
\mathbf{C}_v(n + 1)_{ij} = \mathbf{C}_v(n)_{ij} + (\eta^2 \mu_2 - \eta r_{mod}) (\mathbf{C}_v(n)_{ii} + \mathbf{C}_v(n)_{jj}) + (2\eta^2 \mu_4 - \eta r_{mod}) \sum_{\ell = 1}^{M} (\mathbf{C}_v(n)_{i\ell} + \mathbf{C}_v(n)_{j\ell}) \tag{43}
\]

where \(r_{mod} = [\mathbf{R}_{\mathbf{v},v}]_{ii}\) and \(r_{mod} = [\mathbf{R}_{\mathbf{v},v}]_{ij}\) with \(j \neq i\), as defined in (16).

### E. Some Useful Inequalities

Before concluding this section, let us derive inequalities relating the fourth-order moments \(\mu_i\) and the entries of the covariance matrix \(\mathbf{R}_{\mathbf{v},v}\). For real random variables \(X\) and \(Y\), remember that Hölder’s inequality says \([15]\)
\[
E\{XY\} \leq E\{|X|^{p}\}^{\frac{1}{p}} E\{Y|^{q}\}^{\frac{1}{q}} \tag{44}
\]
where \(p\) and \(q\) are in \([1, \infty)\) with \(\frac{1}{p} + \frac{1}{q} = 1\). As shown hereafter, this inequality yields
\[
\mu_5 \leq \mu_4 \leq \mu_3 \leq \mu_2 \leq \mu_1 \tag{45}
\]

Inequality \((a)\) is obtained by using (44) with \(X = \mathbf{C}_v(n)\), \(Y = \mathbf{C}_v(n)\), \(p = \frac{1}{3}\), and \(q = 4\). In order to prove inequality \((b)\), we first need to observe that Hölder’s inequality yields
\[
E\{X^2Y^2\} \leq E\{|X^3Y\} E\{|XY|^{3}\} \tag{46}
\]
The result \((b)\) directly follows from (46) for \(X = \mathbf{C}_v(n)\) and \(Y = \mathbf{C}_v(n)\). The inequality \((c)\) can be proved using (44) with \(p = \frac{q}{2}, X = \mathbf{C}_v(n)\mathbf{C}_v(n)\) and \(Y = \mathbf{C}_v(n)\mathbf{C}_v(n)\),
\[
E\left\{ \mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n) \right\} \leq E\left\{ \mathbf{C}_v(n)^{1/2}\mathbf{C}_v(n)^{1/2} \right\} E\left\{ \mathbf{C}_v(n)^{1/2}\mathbf{C}_v(n)^{1/2} \right\} \tag{47}
\]
where the equality is due to stationarity of \(\mathbf{C}_v(n)\). Now, for \(p = \frac{q}{2}, X = \mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n)\) and \(Y = \mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n)\) yields
\[
E\left\{ \mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n)\mathbf{C}_v(n) \right\} \leq E\left\{ \mathbf{C}_v(n)^{1/2}\mathbf{C}_v(n)^{1/2} \right\} \times E\left\{ \mathbf{C}_v(n)^{1/2}\mathbf{C}_v(n)^{1/2} \right\} \tag{48}
\]
where the equality is again due to stationarity of \(\mathbf{C}_v(n)\). This last relation proves inequality \((d)\) and completes the proof of (45).

Finally, let us state an inequality involving the main diagonal entry \(r_{mod}\) of the covariance matrix \(\mathbf{R}_{\mathbf{v},v}\). By virtue of Chebyshev’s sum inequality, we have
\[
r_{mod}^2 \leq \mu_3 \tag{49}
\]
3Replace \(X = X^{3/2}Y^{1/2}\) and \(Y = X^{1/2}Y^{3/2}\) in (44) with \(p = q = 2\).
In the next section, we shall use the recursive expressions (42) and (43) of the autocorrelation matrix $C_{\nu}(n_\nu)$, and (45) and (47) to study the steady-state behavior of the Gaussian KLMS algorithm.

IV. GAUSSIAN KLMS ALGORITHM: CONVERGENCE ANALYSIS

We now determine convergence conditions for the Gaussian KLMS algorithm using the analytical model derived in Section III. Let $c_v(n_\nu)$ be the lexicographic representation of $C_{\nu}(n_\nu)$, i.e., the matrix $C_{\nu}(n_\nu)$ is stacked column-wise into a single vector $c_v(n_\nu)$. Consider the family of $(M \times M)$ matrices $H^{ij}$, $1 \leq i, j \leq M$, whose elements are given by (48)–(49), shown at the bottom of the page. Finally, we define the $(M^2 \times M^2)$ matrix $G$ as follows:

$$G = [h^{11} \ h^{12} \ \ldots \ \ h^{1M} \ \ldots \ \ h^{MM}]$$

with $h^{ij}$ the $(M^2 \times 1)$ lexicographic representation of $H^{ij}$. Using these definitions, it can be shown that the lexicographic representation of the recursion (30a) can be written as

$$c_v(n_\nu + 1) = G c_v(n_\nu) + \eta^2 J_{\min} r_{\kappa \kappa}$$

where $r_{\kappa \kappa}$ is the lexicographic representation of $R_{\kappa \kappa}$. We shall use this expression to derive stability conditions for the Gaussian KLMS algorithm. But before closing this subsection, let us remark that matrix $G$ is symmetric. This can be shown from (48)–(49), using $H^{ij}$, and observing that $H^{ij} = H^{ji}$. This implies that $G$ can be diagonalized by a unitary transformation, and all its eigenvalues are real-valued.

A. Convergence Conditions

A necessary and sufficient condition for convergence of $c_v(n_\nu)$ in (51) is that all the eigenvalues of $G$ lie inside the open interval $(-1, 1)$ ([16], Section 5.9). Thus, the stability limit for $\eta$ can be numerically determined for given values of $M$ and $\xi$. In the following we derive a set of analytical sufficient conditions that can also be used for design purposes.

It is well known that all the eigenvalues of $G$ lie inside the union of Gerschgorin disks [17]. Each of these disks is centered at a diagonal element of $G$ and has a radius given by the sum of the absolute values of the remaining elements of the same row. A sufficient condition for stability of (51) is thus given by

$$|G|_{ii} + \sum_{\ell=1, \ell\neq i}^{M^2} |G|_{i\ell} < 1, \quad \text{for } i = 1, \ldots, M^2. \quad (52)$$

Equations (48)–(50) show that the rows of $G$ have only two distinct forms, in the sense that each row of $G$ has the same entries as one of these two distinct rows, up to a permutation. This implies that only two distinct Gerschgorin disks can be defined. Also, all the entries of $G$ are positive except possibly for $G|_{i\ell} = \eta^2 \mu_2 - \eta \tau_{\text{ord}}$ and $G|_{i\ell} = \frac{1}{2}(\eta^2 \mu_4 - \eta \tau_{\text{ord}})$. See Appendix B for proof. Expression (52) thus leads to only two inequalities, defined as follows for $M \geq 3$,

$$\begin{align*}
(1 - 2\eta \tau_{\text{ord}} + \eta^2 \mu_1) + (M - 1) \eta^2 \mu_3 \\
+ 2(M - 1)(M - 2) \eta^2 \mu_4 < 1 \\
(1 - 2\eta \tau_{\text{ord}} + \eta^2 \mu_3) \\
+ 2\eta^2 \mu_2 - \eta \tau_{\text{oc}} + (M - 2) \eta^2 \mu_4 \\
+ 2(M - 2)(2\eta^2 \mu_3 - \eta \tau_{\text{ord}}) \\
+ (M - 2)(M - 3) \eta^2 \mu_5 < 1.
\end{align*} \quad (53a)$$

The study of the limiting values of $\eta$ that satisfy both (53a) and (53b) yields bounds on the maximum step size. The intersection of these bounds yields the sufficient stability limits. In what follows, we present the analysis for $M = 1$ and $M = 2$ are presented in Appendix C.

To determine the bounds imposed by (53a), we write it for $\eta > 0$ as

$$\begin{align*}
2(M - 1) \eta \mu_2 - \eta \tau_{\text{ord}} < 2\tau_{\text{ord}} \\
- \eta \mu_1 - (M - 1) \eta \mu_3 - (M - 1)(M - 2) \eta \mu_4.
\end{align*} \quad (54)$$

Thus, the following two conditions must be satisfied

$$\begin{align*}
2(M - 1)(\eta \mu_2 - \eta \tau_{\text{ord}}) < 2\tau_{\text{ord}} - \eta \mu_1 - (M - 1) \eta \mu_3 - (M - 1)(M - 2) \eta \mu_4 \quad (55a) \\
2(M - 1)(\eta \mu_2 - \eta \tau_{\text{ord}}) > 2\tau_{\text{ord}} + \eta \mu_1 + (M - 1) \eta \mu_3 + (M - 1)(M - 2) \eta \mu_4. \quad (55b)
\end{align*}$$

if $i = j$:

$$\begin{align*}
H_{ii}^{ij} &= 1 - 2\eta r_{\text{adv}} + \eta^2 \mu_1, \\
H_{ii}^{ij} &= \eta^2 \mu_3, \\
H_{ii}^{ij} &= \eta^2 \mu_2 - \eta r_{\text{ord}}, \\
H_{ii}^{ij} &= \eta^2 \mu_4,
\end{align*} \quad (48)$$

if $i \neq j$:

$$\begin{align*}
H_{ii}^{ij} &= 1 - 2\eta r_{\text{adv}} + 2\eta^2 \mu_3, \\
H_{ii}^{ij} &= \eta^2 \mu_4, \\
H_{ii}^{ij} &= \frac{1}{2}(1 - 2\eta r_{\text{adv}} + 2\eta^2 \mu_3), \\
H_{ii}^{ij} &= \eta^2 \mu_5,
\end{align*} \quad (49)$$

$$\begin{align*}
H_{ii}^{ij} &= \frac{1}{2}(\eta^2 \mu_4 - \eta r_{\text{ord}}), \\
H_{ii}^{ij} &= \frac{1}{2}(2\eta^2 \mu_3 - \eta r_{\text{ord}}), \\
H_{ii}^{ij} &= \frac{1}{2}(2\eta^2 \mu_4 - \eta r_{\text{ord}}), \\
H_{ii}^{ij} &= \frac{1}{2}(2\eta^2 \mu_5 - \eta r_{\text{ord}}),
\end{align*}$$

otherwise.
Condition (55a) yields (56), shown at the bottom of the page, because the numerator and denominator of $\eta_1$ are positive. On the other hand, (55b) leads us to a condition of the form $\theta_2 \eta < \theta_1$ with
\[
\begin{align*}
\theta_1 &= 2\tau_{\text{re,d}} - 2(M - 1)\tau_{\text{rd}} \\
\theta_2 &= \mu_1 - 2(M - 1)\mu_2 \\
&\quad + (M - 1)\mu_3 + (M - 1)(M - 2)\mu_4.
\end{align*}
\] (57)
(58)

Hereafter, we shall denote by $\eta_2 := \frac{\theta_2}{\theta_1}$ the resulting bound on $\eta$. Then, solving (55b) yields the four possible cases shown in Table I, depending on the signs of $\theta_1$ and $\theta_2$. In case (ii), note that there exists no $\eta > 0$ that satisfies (55b) when $\theta_1 < 0$ and $\theta_2 > 0$. This situation arises because condition (53a) is a sufficient condition imposing that the Gerschgorin disks defined by (52) are completely inside the unit circle in the z-plane. This condition is obviously not necessary for having all the eigenvalues of $G$ inside the unit circle. The lower bound in case (iv) of Table I is also due to this excessive restriction, as the algorithm is certainly stable for $\eta = 0$. It is kept here for completeness, but it should be disregarded in practice.

Combining the possible solutions of (55a) and (55b) and disregarding lower bounds yields the following stability bounds for $\eta$
\[
\begin{align*}
\eta &< \min\{\eta_1, \eta_2\}, \quad \text{for case (i) in Table I} \\
\eta &< \eta_1, \quad \text{for cases (ii) and (iv) in Table I}.
\end{align*}
\] (59)

If case (iii) happens, which should be tested right away, the stability conditions should be determined numerically through the eigenvalues of $G$.

Having determined limits for (53a) such that $\eta > 0$, and assuming that case (iii) in Table I did not happen, we proceed to determine the extra restrictions imposed on $\eta$ by (53b). First, we multiply (53b) by $(M - 1)$ and divide by $\eta$, rewriting it as
\[
\begin{align*}
2(M - 1)\mu_3 \eta - \tau_{\text{rd}} \\
&\quad < 2(M - 1)\tau_{\text{rd}} - 2(M - 1)\mu_3 \eta \\
&\quad - (M - 1)(M - 2)\mu_4 \eta \\
&\quad - 2(M - 1)(M - 2)\mu_4 \eta - \tau_{\text{re,d}} \\
&\quad - (M - 1)(M - 2)(M - 3)\eta \mu_5.
\end{align*}
\] (60)

Now, given that (54) has already been replaced, we replace the left-hand side (LHS) of (60) with the right-hand side (RHS) of (54). After rearranging the terms, we have the new condition
\[
\begin{align*}
2(M - 1)(M - 2)\mu_4 \eta - 2(M - 1)(M - 2)\tau_{\text{re,d}} \\
&\quad < 2(M - 2)\tau_{\text{re,d}} + \mu_1 \eta \\
&\quad - (M - 1)\mu_3 \eta - (M - 1)(M - 2)(M - 3)\mu_5 \eta \\
&\quad - (M - 1)(M - 2)(M - 3)(M - 4)\mu_5 \eta \\
&\quad - 4(M - 1)(M - 2)\mu_4 \eta - 2(M - 1)(M - 2)\tau_{\text{re,d}}.
\end{align*}
\] (61)

which leads to the following two conditions:
\[
\begin{align*}
4(M - 1)(M - 2)\mu_4 \eta - 2(M - 1)(M - 2)\tau_{\text{re,d}} \\
&\quad < 2(M - 2)\tau_{\text{re,d}} + \mu_1 \eta \\
&\quad - (M - 1)\mu_3 \eta - (M - 1)(M - 2)(M - 3)\mu_5 \eta \\
&\quad - (M - 1)(M - 2)(M - 3)(M - 4)\mu_5 \eta \\
&\quad - 4(M - 1)(M - 2)\mu_4 \eta - 2(M - 1)(M - 2)\tau_{\text{re,d}}.
\end{align*}
\] (62a)
(62b)

On the one hand, (62a) yields two different conditions defined by
\[
\begin{align*}
\eta < \frac{2(M - 2)(\tau_{\text{re,d}} + (M - 1)\tau_{\text{rd}})}{\theta_3}, \quad \text{if } \theta_3 > 0 \\
\eta > 0, \quad \text{if } \theta_3 < 0
\end{align*}
\] (63)

where
\[
\begin{align*}
\theta_3 &= -\mu_1 + (M - 1)\mu_3 + 4(M - 1)(M - 2)\mu_4 \\
&\quad + (M - 1)(M - 2)(M - 3)\mu_5.
\end{align*}
\] (64)

On the other hand, (62b) leads us to a condition of the form $\theta_4 \eta < (M - 2)\theta_1$, where $\theta_1$ has already been defined in (57) and
\[
\begin{align*}
\theta_4 &= -\mu_1 + (M - 1)\mu_3 - 4(M - 1)(M - 2)\mu_4 \\
&\quad + (M - 1)(M - 2)(M - 3)\mu_5.
\end{align*}
\] (65)

Solving (62b) thus leaves four possible cases to consider, shown in Table II, where $\eta_4 := \frac{(M - 2)\theta_1}{\theta_4}$.

\[
\begin{align*}
\eta < \frac{2\tau_{\text{re,d}} + 2(M - 1)\tau_{\text{rd}}}{\mu_1 + 2(M - 1)\mu_2 + (M - 1)\mu_3 + (M - 1)(M - 2)\mu_4} := \eta_1
\end{align*}
\] (56)
Combing the possible solutions of (62a) and (62b), and again disregarding lower bounds, yields the following stability bounds on \( \eta \) [see (66) at the bottom of the page].

Finally, except for cases (iii) in Table I and (vii) in Table II, which should be tested right away, the sufficient stability conditions will be given by the intersection of (59) and (66). These conditions can be slightly simplified by observing that \( \eta_2 \leq \eta_1 \), as can be easily proved using (45).

In the next section, we shall derive the expression of the weight-error correlation matrix \( \mathbf{C}(\eta) \) in steady state. This will allow us to calculate the MSE and the excess MSE.

### B. Steady-State Behavior

The closed-form solution of (51) can be written as [16]

\[
\mathbf{c}_e(n) - \mathbf{G}^n [\mathbf{c}_e(1) - \mathbf{c}_e(\infty)] + \mathbf{c}_e(\infty)
\]

where \( \mathbf{c}_e(\infty) \) denotes the vector \( \mathbf{c}_e(n) \) in steady state, and is given by

\[
\mathbf{c}_e(\infty) = \eta^2 J_{\text{min}} (\mathbf{I} - \mathbf{G})^{-1} \mathbf{r}_{\text{sn}}.
\]

Assuming convergence, we define the time for convergence as the number of iterations required for (67) to reach

\[
|\mathbf{c}_e(n) - \mathbf{c}_e(\infty)| \leq \epsilon
\]

where \( \epsilon \) is a design parameter to be chosen by the user.

At this point, it is important to note that \( \mathbf{c}_e(\infty) \) is unique if the system under consideration is stable. Indeed, the matrix \( (\mathbf{G} - \mathbf{I}) \) has only nonzero eigenvalues because it satisfies conditions (53a)–(53b), and can thus be inverted. Let \( \mathbf{C}_e(\infty) \) be the matrix whose lexicographic representation is \( \mathbf{c}_e(\infty) \). It is also unique and satisfies the following expression derived from (30a) for \( n \rightarrow \infty \):

\[
\mathbf{R}_{\text{sn}} \mathbf{C}_e(\infty) + \mathbf{C}_e(\infty) \mathbf{R}_{\text{sn}} - \eta \mathbf{T}(\infty) - \eta \mathbf{R}_{\text{sn}} J_{\text{min}} = 0.
\]

In (70), \( \mathbf{T}(\infty) \) denotes the matrix \( \mathbf{T}(\eta) \) in steady state. It is interesting to note that (68) is the lexicographic counterpart of (70), as (51) is the lexicographic representation of (30a).

From (23), the steady-state MSE is given by

\[
J_{\text{ss}}(\infty) = J_{\text{min}} + \text{trace}[\mathbf{R}_{\text{ss}} \mathbf{C}_e(\infty)]
\]

where \( \text{trace}[\mathbf{R}_{\text{ss}} \mathbf{C}_e(\infty)] \) is the steady-state excess MSE, denoted by \( J_{\text{ss}}(\infty) \). To evaluate these quantities, we shall now compute the entries of \( \mathbf{C}_e(\infty) \). In order to achieve this, we first justify that \( \mathbf{C}_e(\infty) \) is a matrix with the same structural properties as \( \mathbf{R}_{\text{ss}} \), namely, all its main diagonal entries are equal to each other, having a value denoted by \( C_{e,i} \), and all its off-diagonal entries are also equal to each other, having a value denoted by \( c_{e,i,j} \). Then, we determine \( c_{e,i} \) and \( c_{e,i,j} \) so that \( \mathbf{C}_e(\infty) \) is the solution to (70), which we know is unique. It is straightforward to see that the LHS of (70) is also a matrix with the same structural properties as \( \mathbf{R}_{\text{ss}} \), and as it is equal to \( \eta \mathbf{R}_{\text{ss}} J_{\text{min}} \).

One way for \( \mathbf{C}_e(\infty) \) to have the proposed structure would be that \( \mathbf{R}_{\text{ss}} \mathbf{C}_e(\infty), \mathbf{C}_e(\infty) \mathbf{R}_{\text{ss}} \) and \( \mathbf{T}(\infty) \) all have the same structure. It is straightforward to show that matrices \( \mathbf{R}_{\text{ss}} \mathbf{C}_e(\infty) \) and \( \mathbf{C}_e(\infty) \mathbf{R}_{\text{ss}} \) have this structure if, and only if, \( \mathbf{C}_e(\infty) \) has that structure. If this is the case, a direct consequence, through (70), would be that \( \mathbf{T}(\infty) \) also has this structure. Using \( c_{e,i} \) and \( c_{e,i,j} \), respectively, for the main-diagonal and off-diagonal entries of \( \mathbf{C}_e(\infty) \) in (40) and (41) yields, for \( M > 1 \) [see (72) at the bottom of the page]. For \( M = 1, \mathbf{T}(\infty) = \mu_1 c_{e,i} \). Writing, for ease of notation, \( \mathbf{T}_i(\infty) = t_i c_{e,i} + t_2 c_{e,ij} \) and \( \mathbf{T}_j(\infty) = t_j c_{e,ij} + t_4 c_{e,ij} \) for all \( i \neq j \), solving (70) for \( c_{e,i} \) and \( c_{e,ij} \) yields (73)–(74), shown at the bottom of the page, with \( r_{ei} \) and \( r_{eij} \) defined in (16).

It can be verified that using (72)–(74) in the LHS of (70) yields the correct RHS. Also, we know that this solution is unique since \( \mathbf{I} - \mathbf{G} \) can be inverted. Then, going back to (71), we obtain the following desired result

\[
J_{\text{ss}}(\infty) = J_{\text{min}} + J_{\text{ss}}(\infty)
\]

\[
= J_{\text{min}} + M(r_{\text{md}} c_{\text{md}} + (M - 1)r_{\text{od}} c_{\text{od}}).
\]

(75)

Based on the results presented in the previous sections, we shall now propose design guidelines to set the parameters of

\[
J_{\text{ss}}(\infty) = J_{\text{min}} + \text{trace}[\mathbf{R}_{\text{ss}} \mathbf{C}_e(\infty)]
\]

(71)
the Gaussian KLMS algorithm in order to achieve a prescribed performance.

V. DESIGN GUIDELINES

The analysis results are now used to establish design guidelines. Without limiting generality, the coherence-based sparsification rule (3) is considered hereafter to design the dictionaries. It is however implicit that any other existing technique to select the kernel functions \( \kappa(\cdot, \cdot) \) could be used. These design guidelines assume the availability of multiple realizations of the sequences \( u(n) \) and \( d(n) \).

Suppose a design goal is to obtain an MSE which is less than a specified value \( J_{\text{max}} \). The following procedure could be applied.

1) Set a coherence threshold \( \epsilon_0 \) and define a set of kernel bandwidths \( \xi \) to be tested. For the following design examples, a set of values for \( \xi \) was chosen as equally spaced points in \((0, 1)\). Then, \( \epsilon_0 \) was chosen to yield reasonable values of \( M \) for the chosen set of \( \xi \) values. The value of \( M \) is determined for each pair \( (\xi, \epsilon_0) \) by training the dictionary with the input signal \( u(n) \) until its size stabilizes. The training is repeated several times. A value \( M_i \) is determined for the \( i \)th realization. The value of \( M \) associated with the pair \( (\xi, \epsilon_0) \) is the average of the dictionary sizes for all realizations, rounded to the nearest integer. This is the value of \( M \) to be used in the theoretical model.

2) Using the system input \( u(n) \), determine the desired output \( d(n) \) and estimate \( E\{d^2(n)\} \) and \( p_{nd} \) over several runs. In a practical setting \( d(n) \) should be measured at the unknown system output.

3) From (10), determine the minimum MSE \( J_{\text{in, min}} \) for each set of parameters \((M, \xi)\). If none of the pairs \((M, \xi)\) leads to \( J_{\text{max}} > J_{\text{in, min}} \), return to step 1 to choose a new set of parameters \( \epsilon_0 \) and \( \xi \).

4) Determine \( \eta_{\text{max}} \) using rules (59) and (66) for each value of \( \xi \) and choose the largest possible value of \( \eta \) so that \( \eta < \eta_{\text{max}} \) and \( J_{\text{in, min}}(\infty) < J_{\text{max}} \) (note that \( J_{\text{in, min}}(\infty) \) is computed using (75)).

5) Given \( M, \eta_{\text{max}} \), \( p_{nd} \), \( J_{\text{in, min}} \), and \( \eta \) determine \( G \) using (50) and then \( e_i(n) \) from (51).

6) Determine \( C_\epsilon(\infty) \) from (73)–(74), and then \( J_{\text{max}}(\infty) \); from (75).

7) Choose \( \epsilon \) in criterion (69) and find \( n_e \) from simulations using the parameters determined in steps 1 to 8. Repeat steps 1 to 9 for all kernel parameters in the chosen set and use the one that leads to the smallest \( n_e \).

VI. SIMULATION RESULTS

This section presents examples to illustrate the proposed design procedure and to verify the theoretical results. The simulation conditions which are common to all examples are described in Example 1. Only the changes to these conditions are described in the remaining examples.

A. Example 1

Consider the problem studied in [18] and [19], for which

\[
\begin{align*}
\{ y(n) &= \frac{y(n-1)}{1 + y^2(n-1)} + u^2(n-1) \\
 d(n) &= y(n) + z(n) \}
\end{align*}
\]

(76)

where the output signal \( y(n) \) was corrupted by a zero-mean white Gaussian noise \( z(n) \) with variance \( \sigma^2_z = 10^{-5} \). The input sequence \( u(n) \) is zero-mean i.i.d. Gaussian with standard deviation \( \sigma_u = 0.15 \).

The proposed method was tested with a maximum MSE \( J_{\text{max}} = -21.5 \) dB, a coherence level \( \epsilon_0 = 10^{-2} \) and a set of kernel bandwidths \( \xi \in \{0.01, 0.025, 0.05\} \). For each value of \( \xi \), 500 dictionary dimensions \( M_i \), \( i = 1, \ldots, 500 \), were determined using 500 realizations of the input process. The length of each realization was 500 samples. Each \( M_i \) was determined as the minimum dictionary length required to achieve the coherence level \( \epsilon_0 \). The value \( M(\xi) \) was determined as the average of all \( M_i \), rounded to the nearest integer. The values of \( J_{\text{in, min}}(\xi) \) were calculated from (10) for each pair \( (\xi, M) \). To this end, second-order moments \( p_{nd} \) and \( E\{d^2(n)\} \) were estimated by averaging over 500 runs.

Before searching for a value of \( \eta \), such that \( J_{\text{in, min}}(\infty) < J_{\text{max}} \), we have to define \( \eta_{\text{max}} \). There are two possibilities for defining \( \eta_{\text{max}} \): 1) use the Gerschgorin disk analysis yielding a value of \( \eta_{\text{max}} \) denoted as \( \eta_{\text{max}}^{GD} \) and compute and test the eigenvalues of \( G \) yielding a value of \( \eta_{\text{max}} \) denoted as \( \eta_{\text{max}}^{ED} \). Table III shows that, as expected, the condition imposed by the Gerschgorin disks is more restrictive than that imposed by the eigenvalues of \( G \).

However, note that choosing \( \eta \) from \( \eta_{\text{max}} = \eta_{\text{max}}^{GD} \) is simpler and usually yields good design results.

Table IV presents the obtained results for the chosen values of \( \xi \). For each pair \( (\xi, M) \), the step-size \( \eta \) was chosen so that the algorithm was stable (\( \eta \) less than \( \eta_{\text{max}} \)) and \( J_{\text{max}}(\infty) < -21.5 \) dB. The values of \( J_{\text{in, min}}(\infty) \) and \( J_{\text{in, max}}(\infty) \) were determined from (75) and \( n_e \) was obtained from (69) for \( \epsilon = 10^{-3} \). Note that \( J_{\text{in, min}}(\xi) < -21.5 \) dB in all cases. It clearly appears that \( (\xi, M, \eta) = (1.3, 0.01, 0.152) \) is a good design choice, as it satisfies all design conditions after only \( n_e = 914 \) iterations. Note that the value of \( \eta \) chosen is about 1/10 of \( \eta_{\text{max}} \). This is due to the small value of \( J_{\text{in, min}}(\xi) \) imposed by the design. This is the same phenomenon that happens when designing the regular LMS algorithm for practical specifications in linear estimation problems.

For each simulation, the order \( M \) of the dictionary remained fixed. It was initialized for each realization by generating input
vectors in $\mathcal{U}$ and filling the $M$ positions with vectors that satisfy the desired coherence level. Thus, the initial dictionary is different for each realization. During each realization, the dictionary elements were updated at each iteration $n$ so that the least recently added element is replaced with $u[n - 1]$.

Figs. 2 and 3 illustrate the accuracy of the analytical model for the four cases presented in Table IV. Fig. 2 shows an excellent agreement between Monte Carlo simulations, averaged over 500 runs, and the theoretical predictions made by using (23) and (67). There is again an excellent agreement between theory and simulations.
TABLE V
STABILITY RESULTS FOR EXAMPLE 2

<table>
<thead>
<tr>
<th>ξ</th>
<th>M</th>
<th>J_{\text{min}} [dB]</th>
<th>J_{\text{max}} [\infty] [dB]</th>
<th>J_{\text{ex}} [\infty] [dB]</th>
<th>n_e</th>
</tr>
</thead>
<tbody>
<tr>
<td>0.05</td>
<td>7</td>
<td>2.33</td>
<td>-20.32</td>
<td>-20.27</td>
<td>1544</td>
</tr>
<tr>
<td>0.065</td>
<td>4</td>
<td>2.50</td>
<td>-20.41</td>
<td>-34.63</td>
<td>314</td>
</tr>
<tr>
<td>0.075</td>
<td>3</td>
<td>2.40</td>
<td>-20.68</td>
<td>-33.33</td>
<td>184</td>
</tr>
<tr>
<td>0.125</td>
<td>2</td>
<td>2.34</td>
<td>-20.54</td>
<td>-38.39</td>
<td>502</td>
</tr>
</tbody>
</table>

TABLE VI
SUMMARY OF SIMULATION RESULTS FOR EXAMPLE 2

<table>
<thead>
<tr>
<th>ξ</th>
<th>M</th>
<th>\eta</th>
<th>J_{\text{min}} [dB]</th>
<th>J_{\text{max}} [\infty] [dB]</th>
<th>J_{\text{ex}} [\infty] [dB]</th>
<th>n_e</th>
</tr>
</thead>
<tbody>
<tr>
<td>0.05</td>
<td>7</td>
<td>0.072</td>
<td>-20.32</td>
<td>-20.27</td>
<td>-39.73</td>
<td>1544</td>
</tr>
<tr>
<td>0.065</td>
<td>4</td>
<td>0.239</td>
<td>-20.41</td>
<td>-34.63</td>
<td>314</td>
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<td>0.075</td>
<td>3</td>
<td>0.340</td>
<td>-20.68</td>
<td>-33.33</td>
<td>184</td>
<td></td>
</tr>
<tr>
<td>0.125</td>
<td>2</td>
<td>0.073</td>
<td>-20.54</td>
<td>-38.39</td>
<td>502</td>
<td></td>
</tr>
</tbody>
</table>

B. Example 2

As a second design example, consider the nonlinear dynamic system identification problem studied in [20] where the input signal was a sequence of statistically independent vectors

\[ \mathbf{u}(n) = [u_1(n) \ u_2(n)]^T \] (77)

with correlated samples satisfying \( u_1(n) = 0.5u_0(n) + \eta_u(n) \). The second component of \( \mathbf{u}(n) \) in (77) is an i.i.d. Gaussian noise sequence with variance \( \sigma^2_{u_2} = 0.0156 \) and \( \eta_u(n) \) is a white Gaussian noise with variance \( \sigma^2_{u_1} = 0.0156 \). Consider a linear system with memory defined by

\[ y(n) = \mathbf{a}^T \mathbf{u}(n) - 0.2 y(n-1) + 0.35 y(n-2) \] (78)

where \( \mathbf{a} = [1 \ 0.5]^T \) and a nonlinear Wiener function

\[
\varphi(y(n)) = \begin{cases} 
\frac{y(n)}{3 (0.1 + 0.9 y^2(n))^{1/2}} & \text{for } y(n) \geq 0 \\
-\frac{y^2(n)[1 - \exp(0.7y(n))]^2}{3} & \text{for } y(n) < 0 
\end{cases}
\]

\[ d(n) = \varphi(y(n)) + \epsilon(n) \] (79)

(80)

where \( d(n) \) is the output signal, corrupted by a zero-mean white Gaussian noise \( \epsilon(n) \) with variance \( \sigma^2 = 10^{-6} \). The initial condition \( y(1) = 0 \) was considered in this example.

The proposed method was tested with a maximum MSE \( J_{\text{max}} = -20.25 \) dB, a coherence level \( \epsilon_\eta = 10^{-4} \) and a set of kernel bandwidths \( \xi \in \{0.05, 0.065, 0.075, 0.125\} \). The other simulation conditions were similar to the first example.

Table V shows the estimated values of \( \eta_{\text{max}} \) obtained with the Gershgorin disk analysis and from the eigenvalues of \( \mathbf{G} \). The expression without solution (w.s.) indicates that the intersection of solutions provided by (59) and (66) is empty. In this case, we need to use the limit obtained from the eigenvalues of \( \mathbf{G} \), i.e., \( \eta_{\text{max}} = \eta_{\text{max}}^{\text{w.s.}} \).

Table VI presents the obtained results for the chosen values of \( \xi \). For each pair \( (\xi, M) \), the step-size \( \eta \) was chosen in order to ensure algorithm stability (\( \eta \) less than \( \eta_{\text{max}} \)) and \( J_{\text{ex}} [\infty] < -20.25 \) dB. The values of \( J_{\text{max}} [\infty] \) and \( J_{\text{ex}} [\infty] \) were determined from (75) and \( n_e \) was obtained from (69) for \( \epsilon = 10^{-5} \). Note that \( J_{\text{ex}} [\xi] < -20.25 \) dB in all cases. It clearly appears that \( (M, \xi, \eta) = (3, 0.075, 0.340) \) is a good design choice as it satisfies all design conditions after \( n_e = 184 \) iterations.
Tables VII and VIII present the obtained results for the chosen values of $\xi$. For each pair $(\xi, M)$, the step-size $\eta$ was chosen to ensure stability ($\eta < \eta_{\text{max}}$) and $J_{\text{mex}}(\infty) < -19.75$ dB. The values of $J_{\text{mex}}(\infty)$ and $J_{\text{mex}}'(\infty)$ were determined from (75) and (76), respectively, and $\xi_3$ was obtained from (69) for $\epsilon = 10^{-4}$. Note that $J_{\text{mex}}(\xi) < -19.75$ dB in all cases. It clearly appears that $(M, \xi, \eta) = (7, 0.2, 0.149)$ is a good design choice, as it satisfies the design conditions after $n_t = 382$ iterations.

Figs. 4 and 5 illustrate the accuracy of the analytical model for the four cases presented in Table VI. The agreement between theory and simulations is excellent as in the first example.

C. Example 3

As a third design example, consider the fluid-flow control problem studied in [21] and [22] whose input signal was a sequence of statistically independent vectors

$$u(n) = [u_1(n), u_2(n)]^T$$

with correlated samples satisfying $u_1(n) = 0.5 u_2(n) + \eta_u(n)$, where $u_2(n)$ is a white Gaussian noise sequence with variance $\sigma^2_n = 0.0625^2$ and $\eta_u(n)$ is a white Gaussian so that $\eta_1(n)$ has variance $\sigma_\eta^2 = 0.0625$. Consider a linear system with memory defined by

$$y(n) = a^T u(n) + 1.4138 y(n - 1) - 0.6065 y(n - 2)$$

with $a = [0.1644, 0.0883]^T$ and a nonlinear Wiener function

$$d(n) = \frac{0.3163 y(n)}{0.10 + 0.90 y^2(n)} + z(n)$$

where the output signal $d(n)$ was corrupted by a zero-mean white Gaussian noise $z(n)$ with variance $\sigma^2_z = 10^{-4}$.

The design parameters were set to $J_{\text{max}} = -19.75$ dB (maximum MSE), $\xi_3 = 10^{-1}$ (coherence level) and $\xi \in \{0.15, 0.20, 0.25, 0.30\}$ (set of possible kernel bandwidths).

Table VII shows the values of $\eta_{\text{max}}$ obtained by testing the eigenvalues of $G$. Note that the Gershgorin disk conditions were too strict for all cases making the computation of $\eta_{\text{max}}$ impossible.

Table VIII presents the obtained results for the chosen values of $\xi$. For each pair $(\xi, M)$, the step-size $\eta$ was chosen to ensure stability ($\eta < \eta_{\text{max}}$) and $J_{\text{mex}}(\infty) < -19.75$ dB. The values of $J_{\text{mex}}(\infty)$ and $J_{\text{mex}}'(\infty)$ were determined from (75) and (76), respectively, and $\xi_3$ was obtained from (69) for $\epsilon = 10^{-4}$. Note that $J_{\text{mex}}(\xi) < -19.75$ dB in all cases. It clearly appears that $(M, \xi, \eta) = (7, 0.2, 0.149)$ is a good design choice, as it satisfies the design conditions after $n_t = 382$ iterations.

Figs. 6 and 7 illustrate the accuracy of the analytical model for the four cases presented in Table VIII. Again, the results are very promising for this example.
VII. CONCLUSION

This paper studied the stochastic behavior of the Gaussian KLMS adaptive algorithm for Gaussian inputs and nonlinearities which preserve the stationarity of their inputs. The study resulted in analytical models that predict the behavior of the algorithm as a function of the design parameters. In particular, the new models clarify the joint contribution of the kernel bandwidth and the step-size for the algorithm performance, both during the transient adaptation phase and in steady state. The algorithm convergence was studied and analytical expressions were derived which provide sufficient...
stability conditions. Design guidelines were proposed using the theoretical model. These guidelines were applied to different nonlinear system identification problems. Simulation results have illustrated the accuracy of the theoretical models and their usefulness for design purposes. The extension of this analysis to the case of time-varying systems will be the subject of future investigations.

APPENDIX A

POSITIVE-DEFINITENESS OF $R_{xx}$

To prove that $R_{xx}$ is positive definite we prove that all its eigenvalues are positive [17]. Let $\lambda_i \in \mathbb{H}$ and $\nu_i$, $i = 1, \ldots, M$, be, respectively, the $i$th eigenvalue and the corresponding eigenvector of $R_{xx}$. Hence,

$$R_{xx} \nu_i = \lambda_i \nu_i, \quad i = 1, \ldots, M.$$  \hspace{1cm} (84)

From (16), we can write $R_{xx}$ as

$$R_{xx} = (r_{md} - r_{od}) I + r_{od} 1 1^{\top}$$  \hspace{1cm} (85)

where $1 = [1, 1, \ldots, 1]_{M \times 1}$ and $I$ is the $(M \times M)$ identity matrix. Using (85) in (84) yields

$$[(r_{md} - r_{od}) I + r_{od} 1 1^{\top}] \nu_i = \lambda_i \nu_i, \quad i = 1, \ldots, M.$$  \hspace{1cm} (86)

Noting that $\nu_1 = \beta_1 1$, with $\beta_1 \in \mathbb{R}$ is an eigenvector of $R_{xx}$, we have

$$\beta_1 (r_{md} - r_{od}) I + \beta_1 r_{od} 1 1^{\top} = \lambda_1 \beta_1 1$$  \hspace{1cm} (87)

which yields

$$\lambda_1 = r_{md} + (M - 1) r_{od}$$  \hspace{1cm} (88)

which is positive since $r_{od} > 0$ ($r_{od}$ is real and is the result of a square root in (16)) and $M \geq 1$.

The remaining eigenvectors are orthogonal to $\nu_1$. Thus, $1^{\top} \nu_j = 0$ for all $j \neq 1$ and thus

$$(r_{md} - r_{od}) \nu_j = \lambda_j \nu_j, \quad j = 2, \ldots, M.$$  \hspace{1cm} (89)

which yields

$$\lambda_j = r_{md} - r_{od}, \quad j = 2, \ldots, M.$$  \hspace{1cm} (90)

which are also positive since $r_{md} > r_{od}$. This concludes the proof.

APPENDIX B

SIGN ANALYSIS OF MATRIX $G$ ENTRIES

We shall now analyze the sign of the entries of matrix $G$ defined in (48)–(49). On the one hand, we know that $\eta^2 \mu_3$, $\eta^2 \mu_4$ and $\eta^2 \mu_5$ are strictly positive. On the other hand, $\eta^2 \mu_2 - \eta^2 \mu_3$ and $\frac{1}{2} (2 \eta^2 \mu_4 - \eta^2 \mu_5)$ can be either positive or negative depending on the parameter $\eta$ whose minimum value is $1 - \frac{2 \mu_3}{\mu_1}$. Using (45) and (47), we know that $\mu_4 > r_{md}^2$.

This implies that $1 - 2 \eta r_{md} + \eta^2 \mu_1 > 0$ for all $\eta$. Let us focus now on $\frac{1}{2} (1 - 2 \eta r_{md} + 2 \eta^2 \mu_3)$, whose minimum is equal to $\frac{1}{2} (1 - \frac{2 \mu_3}{\mu_1})$. Similarly as above, we know that $\mu_3 > r_{md}^2$, which means that $\frac{1}{2} (1 - 2 \eta r_{md} + 2 \eta^2 \mu_3) > 0$ for all $\eta$.

As a conclusion, all the diagonal entries of the matrix $G$ are strictly positive, which greatly simplifies the analysis of (52) for stability.

APPENDIX C

STABILITY ANALYSIS IN THE CASES $M = 1$ AND $M = 2$

Let us first derive the condition for stability of the system (51) in the case $M = 1$. The matrix $G$ reduces to the entry $1 - 2 \eta r_{md} + \eta^2 \mu_1$. This directly implies that the system (51) is stable if

$$\eta < \frac{2 r_{md}}{\mu_1}.$$  \hspace{1cm} (91)

Consider now the case $M = 2$. Expression (52) leads us to the following two inequalities:

$$(1 - 2 \eta r_{md} + \eta^2 \mu_1) + \eta^2 \mu_3 + 2 |\eta^2 \mu_2 - \eta r_{md}| < 1,$$  \hspace{1cm} (92a)

$$(1 - 2 \eta r_{md} + \eta^2 \mu_3) + 2 |\eta^2 \mu_2 - \eta r_{md}| < 1.$$  \hspace{1cm} (92b)

We observe that the LHS of (92a) is larger than the LHS of (92b) for all $\eta$ because $\mu_1 \geq \mu_3$, as shown by (45). Proceeding the same way as we did for $M > 3$ in Section IV-A, we conclude that (51) is stable if the following conditions are satisfied:

$$\begin{cases} 
\eta < \min \{\eta_1, \eta_2\}, & \text{for case (i) in Table 1} \\
\eta < \eta_1, & \text{for cases (ii) and (iv) in Table 1}. 
\end{cases}$$  \hspace{1cm} (93)

where

$$\eta_1 = \frac{2 r_{md} + 2 r_{od}}{\mu_1 + 2 \mu_2 + \mu_3},$$  \hspace{1cm} (94a)

$$\eta_2 = \frac{2 r_{md} - 2 r_{od}}{\mu_1 - 2 \mu_2 + \mu_3}.$$  \hspace{1cm} (94b)

REFERENCES


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