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Nonparametric Stochastic Modeling Of Linear Systems With Prescribed Variance Of Several Natural Frequencies

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ABSTRACT: A complete probabilistic model of random positive definite matrices is developed that incorporates constraints on the standard deviations of a set of its eigenvalues. The model is in particular applicable to the representation of the mass and stiffness matrices of random dynamic systems of which certain natural frequencies are observed. The model development is based on the maximization of the entropy under a set of constraints representing the prescribed eigenvalue standard deviations, the mean matrix being given, and the existence of the mean Frobenius norm of the inverse of the random matrix. The efficient simulation of samples of random matrices according to the proposed model is discussed in details. Finally, examples of application validate the above concepts and demonstrate the usefulness of the proposed model.

Keywords: structural dynamics, random systems, random matrices, maximum entropy, probabilistic model

1 INTRODUCTION

The stochastic modeling and simulation of random multi degree of freedom systems has often in the past been accomplished by postulating joint probability density functions of the components of their mass, stiffness, and/or damping matrices or, equivalently, of their corresponding natural frequencies, damping ratios, and mode shapes (e.g. Ghanem and Spanos, 1991, Kleiber et al, 1992, Rivas-Guerra and Mignolet, 2004, Schueller, 1997). However, a different approach has recently been proposed (Soize, 2000, 2001) and validated (see Soize, 2005, for a review) in which the probabilistic model of the mass, stiffness, and/or damping matrices is not assumed but rather determined to maximize the entropy under the constraints (i) that these matrices are positive definite, (ii) that their mean values are prescribed, and (iii) of an overall measure of variation (variance of the norms of the matrices prescribed). This approach has been named nonparametric since no parameter value is selected by the user. It will be further qualified of unconstrained to differentiate it from its constrained extension presented below.

The unconstrained approach is applicable to a broad range of situations in which little, i.e. only a single measure of variability, is known about the randomness of the system. In such cases, the unconstrained approach is particularly advantageous as it

permits to derive, on a rational basis, a full stochastic model of the system. There are however various other situations in which more than a single measure of randomness is known, e.g. when tests have been conducted. Natural frequencies are the most standard information obtained from dynamic tests and thus estimates of the mean and variance of the first few natural frequencies of the system may realistically be available. Since the mean natural frequencies are likely to be close to the values obtained for the design configurations, it is typically the variances which provide the most valuable information about the system randomness. They should thus be incorporated in the stochastic model of the system.

In this light, the goal of the present investigation is to extend the formulation of the unconstrained nonparametric approach to allow for additional constraints on the standard deviations of some of the eigenvalues of the mass, damping, and/or stiffness matrices. This novel approach will be referred to as constrained in the remainder of the paper.

Note that the expected domain of application of these concepts is the low frequency range in which the natural frequencies are distinguishable and observable.

2 ENTROPY MAXIMIZATION AND CONSTRAINTS

It is desired here to simulate realizations of symmetric positive definite random matrices, e.g. the mass and/or stiffness matrix of a linear dynamic system, the properties of which, i.e. eigenvalues, eigenvectors, components, etc., are all random. This requirement necessitates the specification of the joint probability density function of all elements of the matrix. In most practical problems, however, this information is not available - only some moments and/or marginal probability density functions are likely to be available. In the absence of the *exact* distribution, it is then appropriate to ask what are the *desirable* features of this distribution. In this context, note that the design of structural systems is often robust, i.e. that small perturbations in their geometrical and material properties do not alter significantly the probability of failure/fatigue life of the system considered (see Rivas-Guerra and Mignolet, 2004 for a notable counterexample in turbomachinery). It is thus desirable to dispose of a probabilistic model which places particular emphasis on “larger” deviations from the design conditions. Equivalently, this model should have a large value of the entropy as defined by

$$S = - \int_{\Omega} p_{\underline{X}}(\underline{x}) \ln p_{\underline{X}}(\underline{x}) d\underline{x} \quad (1)$$

where \underline{X} denotes the vector of random variables considered of joint probability density function $p_{\underline{X}}(\underline{x})$.

Further, \underline{x} denotes the realized values of \underline{X} , and Ω the domain of support of $p_{\underline{X}}(\underline{x})$.

Consistently with the above discussion, a probabilistic model of uncertain $n \times n$ matrices \underline{A} has been formulated (Soize, 2000, 2001) to maximize the value of the entropy S

$$S = - \int_{\Omega} p_{\underline{A}}(\underline{a}) \ln p_{\underline{A}}(\underline{a}) d\underline{a} \quad (2)$$

given the following physical constraints:

$$\int_{\Omega} p_{\underline{A}}(\underline{a}) d\underline{a} = 1 \quad (3)$$

$$E[\underline{A}] = \int_{\Omega} \underline{a} p_{\underline{A}}(\underline{a}) d\underline{a} = \bar{\underline{A}} \quad (4)$$

and

$$\int_{\Omega} \ln[\det(\underline{a})] p_{\underline{A}}(\underline{a}) d\underline{a} = \nu \text{ finite} \quad (5)$$

where $E[\cdot]$ denotes the operation of mathematical

expectation, $p_{\underline{A}}(\underline{a})$ is the joint probability density function of the elements of \underline{A} , and $\det(\underline{A})$ is its determinant.

The first two of the above constraints correspond to the normalization of the total probability to 1 (Eq. (3)) and the specification of the mean matrix (Eq. (4)). The third one, Eq. (5), implies the existence of the mean squared Frobenius norm of the inverse matrix \underline{A}^{-1} (see Soize (2000,2001) for discussion). To apply this approach to the simulation of random mass, stiffness, and damping matrices of dynamical systems, it is further required to ensure both the symmetry and positive definiteness of every realized matrix \underline{A} . This is achieved by introducing the Cholesky decomposition of \underline{A} , i.e.

$$\underline{A} = \underline{\tilde{L}} \underline{\tilde{L}}^T \quad (6)$$

where $\underline{\tilde{L}}$ is a lower triangular matrix with non-negative diagonal elements and T denotes the operation of matrix transposition. The domain of support Ω of the obtained probability density function is then such that the elements \tilde{L}_{ij} belong to $(-\infty, +\infty)$ for $i \neq j$ and $[0, +\infty)$ for $i = j$, i.e.

$$\Omega = \left\{ \underline{a} = \underline{\tilde{L}} \underline{\tilde{L}}^T ; L_{ij}, i, j = 1, \dots, n : \right. \\ \left. [L_{ij} \in (-\infty, +\infty), i > j] \cap [L_{ii} \in [0, +\infty)] \right\}$$

The maximization of the entropy, Eq. (2), under the constraints of Eqs (3)-(6) yields a closed form expression for the joint probability density function of the elements of the random matrix \underline{A} (Soize, 2000, 2001). Further, this distribution only depends on the single parameter ν so that only a broad knowledge of the matrix uncertainty needs to be known or can be enforced. In some situations however, e.g. when considering insertable turbomachinery blades, tests may have been performed that provide more information on the system variability. In the context of structural dynamics, such tests will often focus on the natural frequencies of the system and will likely result in estimates of the variance of the first few natural frequencies. In such circumstances, it is highly desirable to dispose of a probabilistic model of the corresponding mass and stiffness matrices that accurately accounts for all the available information. If only one natural frequency is observed, its variance as estimated by the tests can serve for the determination of the parameter ν , Eq. (5), corresponding to the mass and/or stiffness matrix model. However, to account for two or more variances, it is necessary to extend the formulation of Eq. (2)-(5) by introducing additional constraints that can reflect the knowledge on the natural fre-

quencies.

To address this extension, consider the generalized eigenvalue problems

$$\underline{\underline{A}}\underline{\underline{\phi}}_i = \lambda_i \underline{\underline{B}}\underline{\underline{\phi}}_i \quad (7)$$

and

$$\overline{\underline{\underline{A}}}\underline{\underline{\phi}}_i = \tilde{\lambda}_i \underline{\underline{B}}\underline{\underline{\phi}}_i \quad (8)$$

where $\underline{\underline{B}}$ is a deterministic symmetric, positive definite matrix. In the ensuing discussions, it will be assumed that the eigenvectors $\underline{\underline{\phi}}_i$ and $\underline{\underline{\phi}}_i$ are normalized with respect to $\underline{\underline{B}}$ so that

$$\underline{\underline{\phi}}_i^T \underline{\underline{B}}\underline{\underline{\phi}}_i = 1 \quad \underline{\underline{\phi}}_i^T \underline{\underline{A}}\underline{\underline{\phi}}_i = \lambda_i \quad (9)$$

and

$$\underline{\underline{\phi}}_i^T \underline{\underline{B}}\underline{\underline{\phi}}_i = 1 \quad \underline{\underline{\phi}}_i^T \overline{\underline{\underline{A}}}\underline{\underline{\phi}}_i = \tilde{\lambda}_i. \quad (10)$$

Constraining the variance of the eigenvalues λ_i directly is unfortunately extremely challenging because of the lack of an exact expression for the natural frequencies of the random matrix $\underline{\underline{A}}$ in terms of its elements. Accordingly, an indirect approach will be selected here which relies on simple constraints that are akin to the second order moments of the eigenvalues. More specifically, it will be assumed that the value of $E\left[\left(\underline{\underline{\phi}}_i^T \underline{\underline{A}}\underline{\underline{\phi}}_i\right)^2\right]$ is specified, i.e.

$$E\left[\left(\underline{\underline{\phi}}_i^T \underline{\underline{A}}\underline{\underline{\phi}}_i\right)^2\right] = \int_{\Omega} \left(\underline{\underline{\phi}}_i^T \underline{\underline{a}}\underline{\underline{\phi}}_i\right)^2 p_{\underline{\underline{A}}}(\underline{\underline{a}}) d\underline{\underline{a}} = s_i^2 \tilde{\lambda}_i^2 \quad (11)$$

where s_i , $i \in I \subseteq [1, n]$, are m known positive constants and $\underline{\underline{\phi}}_i$ are the eigenvectors of the mean matrix $\overline{\underline{\underline{A}}}$ corresponding to the m eigenvalues of which the variance is known. For example, if the variances of the three lowest eigenvalues of $\underline{\underline{A}}$ have been estimated, then $\underline{\underline{\phi}}_i$ will in turn be the eigenvector of the mean matrix $\overline{\underline{\underline{A}}}$ corresponding to its three lowest eigenvalues.

Several comments can be made in regards to the constraints of Eq. (11). First, these conditions involve second order moments, not variances, but this switch is appropriate because the mean values of $\left(\underline{\underline{\phi}}_i^T \underline{\underline{A}}\underline{\underline{\phi}}_i\right)$ are already prescribed by Eq. (4). Next, as discussed above, these conditions do not generally relate exactly to the natural frequencies of the random matrix $\underline{\underline{A}}$, but they do so when its eigenvectors are the same as those of its mean $\overline{\underline{\underline{A}}}$. Finally, it should be noted that the specification of the constraints of Eq. (11) provides freedom in the prob-

abilistic model of the random matrix $\underline{\underline{A}}$ which can be used to match the known variances of the natural frequencies. This last issue will be discussed in more details further.

3 PROBABILISTIC MODEL DERIVATION

Following the discussions of the previous sections, the proposed probabilistic model $p_{\underline{\underline{A}}}(\underline{\underline{a}})$ maximizes the entropy, S of Eq. (2), under the constraints of Eqs (3)-(5) and (11) as well as the symmetry and positive definiteness requirements of Eq. (6). Using Lagrange multipliers μ_0 , $\underline{\underline{\mu}}$, $\lambda - 1$, and $\tilde{\tau}_i$ the constrained maximization of Eq. (2) is reduced to the unconstrained maximization of

$$S^* = S - \mu_0 \int_{\Omega} p_{\underline{\underline{A}}}(\underline{\underline{a}}) d\underline{\underline{a}} - \text{tr} \left\{ \underline{\underline{\mu}}^T \int_{\Omega} \underline{\underline{a}} p_{\underline{\underline{A}}}(\underline{\underline{a}}) d\underline{\underline{a}} \right\} + (\lambda - 1) \int_{\Omega} \ln[\det(\underline{\underline{a}})] p_{\underline{\underline{A}}}(\underline{\underline{a}}) d\underline{\underline{a}} - \sum_{i \in I} \tilde{\tau}_i \int_{\Omega} \left(\underline{\underline{\phi}}_i^T \underline{\underline{a}}\underline{\underline{\phi}}_i\right)^2 p_{\underline{\underline{A}}}(\underline{\underline{a}}) d\underline{\underline{a}} \quad (12)$$

Proceeding next by calculus of variations, it is shown that

$$p_{\underline{\underline{A}}}(\underline{\underline{a}}) = \tilde{C} [\det(\underline{\underline{a}})]^{\lambda-1} \exp \left[-\text{tr}(\underline{\underline{\mu}}^T \underline{\underline{a}}) - \sum_{i \in I} \tilde{\tau}_i \left(\underline{\underline{\phi}}_i^T \underline{\underline{a}}\underline{\underline{\phi}}_i\right)^2 \right] \quad (13)$$

where \tilde{C} is the appropriate constant to satisfy the normalization condition, Eq. (3). It should be noted from Eq. (13) that this probabilistic model is independent of a rotation/change of coordinate system as the vectors $\underline{\underline{\phi}}_i$ are fixed in space.

Before addressing the evaluation of the Lagrange multipliers, it is desired to simplify Eq. (13) and to address the positive definiteness requirement. In regards to simplifications, introduce first the matrix $\overline{\underline{\underline{L}}}$ such that

$$\overline{\underline{\underline{A}}} = \overline{\underline{\underline{L}}}\overline{\underline{\underline{L}}}^T \quad \text{and} \quad \overline{\underline{\underline{L}}}\underline{\underline{\Phi}} = \begin{bmatrix} \tilde{\underline{\underline{\Lambda}}}^{1/2} \\ \underline{\underline{0}} \end{bmatrix} \quad (14a,b)$$

where $\underline{\underline{\Phi}}$ denotes the $n \times m$ matrix formed by the m eigenvectors $\underline{\underline{\phi}}_i$, $i \in I$, and $\tilde{\underline{\underline{\Lambda}}}$ is the corresponding diagonal matrix of eigenvalues. Note that the condition of Eq. (14b) is introduced to simplify the constraints of Eq. (11) as will be shown below. From Eqs (7)-(10), it can be proved that $\overline{\underline{\underline{L}}}$ can be expressed in the partitioned form

$$\underline{\underline{L}} = \left[\underline{\underline{A}} \underline{\underline{\Phi}} \underline{\underline{\Lambda}}^{-1/2} \quad \underline{\underline{D}} \right] \quad (15)$$

where the $n \times (n-m)$ matrix $\underline{\underline{D}}$ is any decomposition, e.g. Cholesky, of

$$\underline{\underline{D}} \underline{\underline{D}}^T = \underline{\underline{A}} - \underline{\underline{A}} \underline{\underline{\Phi}} \underline{\underline{\Lambda}}^{-1} \underline{\underline{\Phi}}^T \underline{\underline{A}}. \quad (16)$$

Pre- and postmultiplying Eq. (16) by $\underline{\underline{\Phi}}^T$ and $\underline{\underline{\Phi}}$, respectively, it is found that $\underline{\underline{D}}^T \underline{\underline{\Phi}} = \underline{\underline{0}}$ as required in Eq. (14b).

Next, express the random matrix $\underline{\underline{A}}$ as

$$\underline{\underline{A}} = \underline{\underline{L}} \underline{\underline{G}} \underline{\underline{L}}^T. \quad (17)$$

Proceeding with this change of random variables, it is found that the probability density function of the elements of $\underline{\underline{G}}$ is

$$p_{\underline{\underline{G}}}(\underline{\underline{g}}) = \bar{C} [\det(\underline{\underline{g}})]^{\lambda-1} \exp \left[-\text{tr}(\underline{\underline{\mu}}^T \underline{\underline{g}}) - \sum_{i=1}^m \tau_i g_{ii}^2 \right] \quad (18)$$

where \bar{C} is a new normalization constant,

$$\underline{\underline{\mu}} = \underline{\underline{L}}^T \underline{\underline{\mu}} \underline{\underline{L}}, \text{ and } \tau_i = \tilde{\tau}_i \tilde{\lambda}_i^2.$$

To guarantee the symmetry and positive definiteness of $\underline{\underline{G}}$, and thus of $\underline{\underline{A}}$, the model of Eq. (18) is reformulated in terms of the elements of the lower triangular matrix $\underline{\underline{H}}$ such that

$$\underline{\underline{G}} = \underline{\underline{H}} \underline{\underline{H}}^T. \quad (19)$$

As demonstrated in Soize (2000), the Jacobian of the transformation is

$$J = 2^n \prod_{l=1}^n h_{ll}^{n-l+1} \quad (20)$$

so that

$$p_{\underline{\underline{H}}}(\underline{\underline{h}}) = C \left\{ \prod_{l=1}^n h_{ll}^{n-l+2\lambda-1} \right\} \times \exp \left[-\text{tr}(\underline{\underline{h}}^T \underline{\underline{\mu}}^T \underline{\underline{h}}) - \sum_{i=1}^m \tau_i \left(\sum_{l=1}^i h_{il}^2 \right)^2 \right] \quad (21)$$

where C is the appropriate normalization constant over the domain

$$\bar{\Omega} = \{h_{ij}, i, j = 1, \dots, n : [h_{ij} \in (-\infty, +\infty), i > j] \cap [h_{ii} \in [0, +\infty)]\}.$$

To evaluate the Lagrange multipliers $\underline{\underline{\mu}}$ and τ_i , it is necessary to first express the constraints of Eqs (4) and (11) in terms of the elements of the random matrix $\underline{\underline{H}}$. Combining Eqs (4), (17), and (19), it is found that

$$E[\underline{\underline{G}}] = \underline{\underline{I}}_n \quad (22)$$

where $\underline{\underline{I}}_n$ denotes the $n \times n$ identity matrix. The simplicity of this condition implies an equally simple

form of the matrix $\underline{\underline{\mu}}$. Specifically, it can be shown that this matrix is diagonal and thus Eq. (21) reduces to

$$p_{\underline{\underline{H}}}(\underline{\underline{h}}) = \prod_{i=1}^m \left\{ C_i h_{ii}^{p(i)} \exp \left[-\mu_{ii} \left(\sum_{l=1}^i h_{il}^2 \right) - \tau_i \left(\sum_{l=1}^i h_{il}^2 \right)^2 \right] \right\} \times \prod_{i=m+1}^n \left\{ C_i h_{ii}^{p(i)} \exp \left[-\mu_{ii} h_{ii}^2 \right] \right\} \times \prod_{i=m+1}^n \prod_{l=1}^{i-1} \left\{ C_{il} \exp \left[-\mu_{ii} h_{il}^2 \right] \right\} \quad (23)$$

where C_i , $i=1, \dots, n$, and C_{il} , $i=m+1, \dots, n$; $l=1, \dots, i-1$, are appropriate normalization constants and

$$p(i) = n - i + 2\lambda - 1. \quad (24)$$

It is concluded from Eq. (23) that:

(i) the elements h_{il} , $i > l$ and $i > m$ are all independent of each other and independent of the other elements h_{il} . Further, they are normally distributed with mean 0 and standard deviation

$$\sigma_{il} = 1 / \sqrt{2\mu_{ii}}.$$

(ii) the elements h_{ii} , $i > m$, are all independent of each other and independent of the other elements h_{il} . Further, they are distributed according to

$$p_{H_{ii}}(h_{ii}) = C_i h_{ii}^{p(i)} \exp \left[-\mu_{ii} h_{ii}^2 \right], h_{ii} \geq 0 \quad (25)$$

where

$$C_i = \frac{2\mu_{ii}^{[p(i)+1]/2}}{\Gamma((p(i)+1)/2)} \quad (26)$$

and $\Gamma(\cdot)$ denotes the Gamma function.

(iii) the elements h_{il} , $l=1, \dots, i$ for a given $i \in [1, m]$ are dependent on each other but independent of the other elements h_{il} . Their joint distribution is

$$p_{H_{il}}(h_{il}) = C_i h_{ii}^{p(i)} \exp \left[-\mu_{ii} \left(\sum_{l=1}^i h_{il}^2 \right) - \tau_i \left(\sum_{l=1}^i h_{il}^2 \right)^2 \right] \quad (27)$$

over the domain

$$\bar{\Omega}_i = \{h_{il}, l=1, \dots, i : [h_{il} \in (-\infty, +\infty), i > l] \cap [h_{ii} \in [0, +\infty)]\}.$$

From the first observation, (i), it is concluded that $E[G_{il}] = 0$ for $i > l$ and $i > m$. Further, the observation (iii) and the symmetry of the distribution of Eq. (27) with respect to the origin imply similarly that $E[G_{il}] = 0$ for $i > l$ and $i \leq m$. It then remains to satisfy the diagonal terms of the mean condition, Eq. (22). For $i > m$,

$$E[G_{ii}] = 1 = \sum_{l=1}^i E[H_{il}^2] = (i-1) \frac{1}{2\mu_{ii}} + E[H_{ii}^2] \quad (28)$$

where $E[H_{ii}^2]$ is obtained by integration of Eq. (27) as $(p(i)+1)/2\mu_{ii}$. Combining this result and Eqs

(26) and (28), it is found that

$$\mu_{ii} = \frac{n+2\lambda-1}{2} \text{ for } i > m. \quad (29)$$

At this point, it only remains to determine the Lagrange multipliers μ_{ii} and τ_i for $i \in [1, m]$. This step is achieved by enforcing the conditions

$$E[G_{ii}] = 1 \quad (30)$$

and Eq. (11) or

$$E[G_{ii}^2] = s_i^2 \quad (31)$$

for $i \leq m$. To this end, it is useful to proceed with the following change of variables that highlights the random element G_{ii}

$$G_{ii} = \sum_{l=1}^i H_{il}^2 \quad (32)$$

$$H_{ii} = \sqrt{G_{ii}} \cos \Theta_{i1} \quad (33)$$

$$H_{i(i-1)} = \sqrt{G_{ii}} \sin \Theta_{i1} \cos \Theta_{i2} \quad (34)$$

$$H_{i(i-2)} = \sqrt{G_{ii}} \sin \Theta_{i1} \sin \Theta_{i2} \cos \Theta_{i3} \quad (35)$$

and so on till

$$H_{i1} = \sqrt{G_{ii}} \sin \Theta_{i1} \sin \Theta_{i2} \dots \sin \Theta_{i(i-1)} \quad (36)$$

where $\Theta_{i1} \in [0, \pi/2]$, $\Theta_{il} \in [0, \pi]$, $l=2, \dots, i-2$, and $\Theta_{i(i-1)} \in [0, 2\pi]$. The Jacobian of this transformation can be found (e.g. see Soize, 2000) to be

$$J' = \frac{(g_{ii})^{i/2-1}}{2} |\sin \theta_{i1}|^{i-2} |\sin \theta_{i2}|^{i-3} \dots |\sin \theta_{i(i-2)}| \quad (37)$$

Then, the joint probability density function of G_{ii} and Θ_{il} , $l=1, \dots, i-1$, is

$$p_{G_{ii} \Theta_{il}}(g_{ii}, \theta_{il}) = \frac{1}{2} C_i g_{ii}^{(n+2\lambda-3)/2} \exp[-\mu_{ii} g_{ii} - \tau_i g_{ii}^2] \times |\sin \theta_{i1}|^{i-2} [\cos \theta_{i1}]^{p(i)} \times |\sin \theta_{i2}|^{i-3} \dots |\sin \theta_{i(i-2)}|. \quad (38)$$

The above expression demonstrates that the random variables G_{ii} and Θ_{il} , $l=1, \dots, i-1$ are all independent of each other and that the joint probability density function of the angles Θ_{il} , $l=1, \dots, i-1$ does not depend on the values of τ_i . This important observation will be used in the next section. Further, the marginal distribution of G_{ii} is

$$p_{G_{ii}}(g_{ii}) = \bar{C}_i g_{ii}^{(n+2\lambda-3)/2} \exp[-\mu_{ii} g_{ii} - \tau_i g_{ii}^2] \quad (39)$$

where $g_{ii} \geq 0$. For each value of i , the two parameters μ_{ii} and τ_i are then finally solved from the constraints given by Eqs (30) and (31).

4 SIMULATION OF RANDOM MATRICES

The simulation of random matrices \underline{A} according to the model derived above is achieved by first generating random matrices \underline{H} according to the joint distributions of Eqs (23), (25), (27), (38) and (39). Once a sample of \underline{H} has been generated, the corresponding matrix \underline{G} is determined from Eq. (19) and, finally, a realization of \underline{A} is obtained from Eq. (17). Thus, the simulation effort reduces to the generation of appropriate samples of \underline{H} . From the observations drawn in the previous section, it is concluded that there are, besides the simulation of the Gaussian variates H_{il} , $i > m$ and $i > l$, three particular issues. These are: (i) the generation of the diagonal elements H_{ii} , $i > m$ according to the distribution of Eq. (25), (ii) the simulation of the ratios H_{il} / G_{ii} , $i \leq m$, described by the angular variables in Eqs (33)-(38), and (iii) the generation of samples of G_{ii} , $i \leq m$, according to Eq. (39). These three issues are addressed in order below.

4.1 Simulation of H_{ii} , $i > m$

The generation of samples of H_{ii} is simplified by considering the variable $Y_{ii} = \mu_{ii} H_{ii}^2$. Proceeding with the change of variables, it is found that the probability density function of Y_{ii} is

$$p_{Y_{ii}}(y_{ii}) = \frac{y_{ii}^{(p(i)-1)/2}}{\Gamma[(p(i)+1)/2]} \exp[-y_{ii}], \quad y_{ii} \geq 0. \quad (40)$$

That is, Y_{ii} is a Gamma distributed random variable for which efficient simulation algorithms exist, e.g. see Devroye (1986). Once a sample of Y_{ii} has been simulated according to the Gamma distribution, the corresponding value of H_{ii} , $i > m$, is found as

$$H_{ii} = \sqrt{\frac{Y_{ii}}{\mu_{ii}}} \quad (41)$$

where μ_{ii} is given by Eq. (29).

4.2 Simulation of $H_{il} / \sqrt{G_{ii}}$, $i \leq m$, $l=1, \dots, i$

A first approach for the simulation of the random variables $H_{il} / \sqrt{G_{ii}}$, $l=1, \dots, i-1$, and $i \leq m$, is to proceed from the generalized spherical coordinates transformation of Eqs (33)-(36) and to generate independent angles Θ_{il} , $l=1, \dots, i-1$, according to the distributions

$$p_{\Theta_{il}}(\theta_{il}) = \tilde{C}_{il} |\sin \theta_{il}|^{i-2} [\cos \theta_{il}]^{p(i)}, \quad \theta_{il} \in \left[0, \frac{\pi}{2}\right] \quad (42)$$

and

$$p_{\Theta_{il}}(\theta_{il}) = \tilde{C}_{il} |\sin \theta_{il}|^{i-l-1} \quad (43)$$

with $\theta_{il} \in [0, \pi]$ for $l=2, \dots, i-2$, and $\theta_{i(i-1)} \in [0, 2\pi)$

and where the coefficients \tilde{C}_{il} are appropriate normalization constants.

Note however that the above distributions are non standard and thus a different, easier approach was selected here. Specifically, it was observed in the previous section that the probability density functions of Eqs (42) and (43) do not depend on τ_i and thus they would be the same for $\tau_i=0$. In this case however, the simulation of the random variables H_{il} is completely similar to the case $i>m$, i.e. the random variables H_{il} , $l \neq i$, are Gaussian variates with mean zero and standard deviation $1/\sqrt{2\mu_{ii}}$ while $H_{ii} = \sqrt{Y_{ii}/\mu_{ii}}$ where Y_{ii} is a Gamma random variable. Note further that the appropriate value of μ_{ii} to be used is the one given by Eq. (29) to insure the consistency with $p(i)$ of Eq. (24).

In view of these comments, the simulation of the terms $H_{il}/\sqrt{G_{ii}}$, $l=1, \dots, i-1$, and $i \leq m$, can efficiently be accomplished as follows:

(a) generate \tilde{H}_{ii} as

$$\tilde{H}_{ii} = \sqrt{\frac{\tilde{Y}_{ii}}{\mu_{ii}}} \quad (44)$$

where μ_{ii} is given by Eq. (29) and \tilde{Y}_{ii} is a Gamma distributed random variable of probability density function given by Eq. (40).

(b) generate \tilde{H}_{il} , $l=1, \dots, i-1$, as zero mean Gaussian random variables with standard deviation $1/\sqrt{2\mu_{ii}}$.

(c) form $\tilde{G}_{ii} = \sum_{l=1}^i \tilde{H}_{il}^2$ and the desired ratios

$$\frac{H_{il}}{\sqrt{G_{ii}}} = \frac{\tilde{H}_{il}}{\sqrt{\tilde{G}_{ii}}} \quad l=1, \dots, i. \quad (45)$$

4.3 Simulation of G_{ii} , $i \leq m$

The last step in the simulation of the random matrix \underline{H} is the generation of the diagonal elements G_{ii} according to their probability density function of Eq. (39). This distribution is non standard and thus a dedicated simulation algorithm by rejection from the

Student's t distribution with 3 degrees-of-freedom was devised (see Devroye (1986) for background). Such an algorithm requires that there exists a constant c_θ such that

$$p_{G_{ii}}(g_{ii}) \leq c_\theta p_{Z_\theta}(g_{ii}) \quad \text{for all } g_{ii} \geq 0 \quad (46)$$

where $p_{G_{ii}}(g_{ii})$ is defined by Eq. (39) and

$p_{Z_\theta}(g_{ii})$ is the probability density function of

$Z_\theta = y_0 + \theta Z$ in which Z is a random variable having the Student's t distribution with 3 degrees-of-freedom. That is, the probability density function of Z is

$$p_Z(z) = \frac{1}{2\sqrt{2}} \left(1 + \frac{z^2}{2}\right)^{-3/2} \quad z \in (-\infty, +\infty) \quad (47)$$

and thus,

$$p_{Z_\theta}(g_{ii}) = \frac{1}{2\sqrt{2}\theta} \left(1 + \frac{(g_{ii} - y_0)^2}{2\theta^2}\right)^{-3/2} \quad (48)$$

for $g_{ii} \in (-\infty, +\infty)$. If values of c_θ , θ , and y_0 can be found for which Eq. (46) holds, then random deviates G_{ii} can be generated as

$$G_{ii} = Z_\theta \quad \text{provided that } c_\theta p_{Z_\theta}(Z_\theta) \bar{U} \leq p_{G_{ii}}(Z_\theta) \quad (49)$$

in which \bar{U} denotes a random variable uniform in the interval $[0,1]$ and independent of Z_θ . When the inequality in Eq. (49) is not satisfied, no sample G_{ii} is generated. Note further that $1/c_\theta$ is the probability that this inequality be satisfied, so that c_θ represents the average number of pairs of samples (Z_θ, \bar{U}) that must be generated per value of G_{ii} to be simulated. Accordingly, it is desired to have a value of c_θ as close as possible to 1.

The Student's t distribution with 3 degrees-of-freedom was selected for the random variable Z first because it can easily be simulated as

$$Z = \sqrt{\frac{2}{U(1-U)}} \left(U - \frac{1}{2}\right) \quad (50)$$

where U is a random number uniformly distributed in $[0,1]$. Another advantage of the Student's t distribution is that it led to values of c_θ that remained reasonably small, i.e. to efficient simulation algorithms, over a broad range of values of the parameters $p(i)$ and τ_i/μ_{ii}^2 .

The approach selected here for the determination of the parameters c_θ , θ , and y_0 is to force the two distributions $p_{G_{ii}}(g_{ii})$ and $p_{Z_\theta}(g_{ii})$ to have their

respective modes at the same value $g_{ii}^* = y_0$ and to have $p_{G_{ii}}(g_{ii}^*) = c_\theta p_{Z_\theta}(g_{ii}^*)$. These 2 conditions will provide the values of c_θ and y_0 in terms of θ . This last parameter will then be selected to minimize c_θ while maintaining the inequality of Eq. (46) over the entire domain $g_{ii} \geq 0$.

Matching the modes of the two distributions $p_{G_{ii}}(g_{ii})$ and $p_{Z_\theta}(g_{ii})$ leads directly to the condition

$$2\tau_i y_0^2 + \mu_{ii} y_0 - q = 0 \quad (51)$$

or

$$y_0 = \frac{1}{4\tau_i} \left(\sqrt{\mu_{ii}^2 + 8q\tau_i} - \mu_{ii} \right) \geq 0 \quad (52)$$

where the notation

$$q = \frac{n + 2\lambda - 3}{2} \quad (53)$$

has been adopted for simplicity.

Next, matching the peak values of $p_{G_{ii}}(g_{ii})$ and $c_\theta p_{Z_\theta}(g_{ii})$ requires that

$$c_\theta = 2\sqrt{2} \theta \bar{C}_i y_0^q \exp[-\mu_{ii} y_0 - \tau_i y_0^2] \quad (54)$$

which indicates that c_θ is proportional to θ while

y_0 is independent of this parameter, see Eq. (52).

An acceptable simulation algorithm is obtained for all values of θ for which the inequality of Eq. (46) is satisfied. However, the most efficient of these algorithms is the one that minimizes the corresponding c_θ , see Eq. (54). In view of the linearity of c_θ with respect to θ , it is concluded that the best algorithm is the one for which θ has the smallest possible value that guarantees the satisfaction of the inequality of Eq. (46) for all values of $g_{ii} \geq 0$. To determine this value of θ , introduce first the function

$$r(w) = \ln \frac{c_\theta p_{Z_\theta}(y_0 + w)}{p_{G_{ii}}(y_0 + w)} \quad (55)$$

and it is desired that $r(w)$ be nonpositive for all $w \geq -y_0$. From the above conditions, it is found that $r(0)=0$ and $r'(0)=0$. Thus, a sufficient condition for $r(w)$ to be nonpositive is that

$$r'(w) \geq 0 \text{ for } w < 0 \text{ and } r'(w) \leq 0 \text{ for } w > 0. \quad (56)$$

After differentiation of Eq. (55) and some algebraic manipulations, it is found that Eq. (56) is equivalent to

$$2\theta^2 \geq \frac{3(y_0 + w)}{(\mu_{ii} + 4\tau_i y_0 + 2\tau_i w)} - w^2 \quad (57)$$

for all $w \geq -y_0$. Accordingly, the smallest value of θ corresponds to the maximum of the right-hand-

side of Eq. (57). In this regards, note from Eqs (51) and (52) that $\mu_{ii} + 4\tau_i y_0 + 2\tau_i w$ is always positive for all $w \geq -y_0$. Then, the right-hand-side of Eq. (57) is monotonically increasing in $w \in [-y_0, 0]$ and thus the maximum must occur for $w > 0$. Further,

$$\begin{aligned} & \max_{w>0} \left\{ \frac{3(y_0 + w)}{(\mu_{ii} + 4\tau_i y_0 + 2\tau_i w)} - w^2 \right\} \quad (58) \\ & \leq \max_{w>0} \left\{ \frac{3(y_0 + w)}{(\mu_{ii} + 4\tau_i y_0)} - w^2 \right\} = \frac{9 + 12y_0(\mu_{ii} + 4\tau_i y_0)}{4(\mu_{ii} + 4\tau_i y_0)^2} \end{aligned}$$

which suggests the value of θ as

$$\theta = \sqrt{\frac{9 + 12y_0(\mu_{ii} + 4\tau_i y_0)}{8(\mu_{ii} + 4\tau_i y_0)^2}}. \quad (59)$$

The above expression is simple but overestimates the maximum of the function on the right-hand-side of Eq. (57). This maximum can be obtained by differentiating and solving a cubic equation in w . Specifically, it is found that

$$\theta = \sqrt{\frac{3(y_0 + w^*)}{2(\mu_{ii} + 4\tau_i y_0 + 2\tau_i w^*)} - \frac{w^{*2}}{2}} \quad (60)$$

where w^* denotes the location of the maximum as given by

$$w^* = \frac{(D - a)}{2\tau_i} \quad (61)$$

where

$$D = \frac{\gamma}{6} + \frac{2a^2}{3\gamma} + \frac{a}{3} \quad (62)$$

$$\gamma^3 = 108b + 8a^3 + 12\sqrt{81b^2 + 12ba^3} \quad (63)$$

and

$$a = \mu_{ii} + 4\tau_i y_0 \quad b = 3\tau_i(\mu_{ii} + 2\tau_i y_0). \quad (64)$$

A comparison of the simulation algorithms based on Eqs (59) and (60)-(64) will be presented in section 6.1.

The simulation of random values G_{ii} according to the probability density function of Eq. (39) then proceeds as follows. For given values of μ_{ii} and τ_i , the mode y_0 is first determined according to Eq. (52), and θ and c_θ / \bar{C}_i are obtained in order from Eqs (59), or (60)-(64), and (54). After this preparation phase, pairs of independent random numbers (U, \bar{U}) uniform in $[0, 1]$ are simulated. From U , a variate Z of the Student's t distribution is obtained from Eq. (50) and the corresponding variable Z_θ is determined as $Z_\theta = y_0 + \theta Z$. It is next necessary to assess if the inequality of Eq. (49) is satisfied. To

this end, the term $[c_\theta / \bar{C}_i] p_{Z_\theta}(Z_\theta) \bar{U}$ is evaluated from the above values of c_θ / \bar{C}_i and \bar{U} and with $p_{Z_\theta}(Z_\theta)$ denoting the value of $p_{Z_\theta}(g_{ii})$, Eq. (48), for $g_{ii} = Z_\theta$. The ratio $p_{G_{ii}}(Z_\theta) / \bar{C}_i$ is similarly obtained from Eq. (39). If

$$[c_\theta / \bar{C}_i] p_{Z_\theta}(Z_\theta) \bar{U} \leq p_{G_{ii}}(Z_\theta) / \bar{C}_i, \quad (65)$$

a value of G_{ii} is obtained as $G_{ii} = Z_\theta$. Otherwise, the value of Z_θ is rejected and no corresponding sample of G_{ii} is generated. Either way, the process is repeated starting with the simulation of a new pair of uniform random variables (U, \bar{U}) until the appropriate number of samples of G_{ii} has been obtained. Note that Eq. (65) is equivalent to the inequality of Eq. (49) but is preferable to it because it does not require the numerical evaluation of the normalization constant \bar{C}_i .

5 IDENTIFICATION OF NONPARAMETRIC MODEL PARAMETERS

5.1 Dispersion Parameter δ

It is desired next to characterize the overall variability of the random matrices \underline{A} around their mean $\bar{\underline{A}}$. To this end, introduce the dispersion parameter δ as

$$\delta^2 = \frac{1}{n} E \left[\left\| \underline{G} - \frac{1}{n} \underline{I} \right\|_F^2 \right] \quad (66)$$

where

$$\left\| \underline{V} \right\|_F^2 = \text{tr} \left[\underline{V} \underline{V}^T \right] \quad (67)$$

is the square of the Frobenius norm of an arbitrary matrix \underline{V} . Expanding, the right-hand-side of Eq. (67) and taking Eq. (22) into account, it is found that

$$\delta^2 = \frac{1}{n} \left\{ E \left[\text{tr} \left(\underline{G} \underline{G}^T \right) \right] - n \right\} \quad (68)$$

with

$$E \left[\text{tr} \left(\underline{G} \underline{G}^T \right) \right] = \sum_{i,j,k,l=1}^n E \left[H_{ik} H_{jk} H_{il} H_{jl} \right] \quad (69)$$

where, by its lower triangular nature, $H_{il} = 0$ for $l > i$. The evaluation of the right-hand-side of Eq. (69) is accomplished by first separating the cases in which $i=j$ and those in which $i \neq j$. Further, it is recalled from Eq. (23) that H_{ik} and H_{jl} are independent if $i \neq j$. Noting finally from Eq. (23) that $E[H_{ik} H_{il}] = 0$ for $k \neq l$, it is found that

$$E \left[\text{tr} \left(\underline{G} \underline{G}^T \right) \right] = \sum_{i=1}^n E \left[G_{ii}^2 \right] - \sum_{i,k=1}^n \left(E \left[H_{ik}^2 \right] \right)^2 + \sum_{k=1}^n \left(E \left[\sum_{i=k}^n H_{ik}^2 \right] \right)^2. \quad (70)$$

To evaluate the above expression, it should further be noted that $E[H_{ik}^2] = E[H_{i(i-1)}^2]$, for $k \neq i$, in view of the symmetry of the variables H_{il} , $l > i$, in Eq. (23) and further that these expectations are the same if $i > m$ or $i \leq m$ in view of Eq. (30). Moreover, from Eqs (30), (33), and (34), it is found that

$$E \left[H_{ii}^2 \right] = \frac{n + 2\lambda - i}{n + 2\lambda - 1} \quad (71)$$

and

$$E \left[H_{ii}^2 \right] = \frac{1}{n + 2\lambda - 1}. \quad (72)$$

Introducing Eq. (71) and (72) into (70) and combining the above results yields

$$\delta^2 = \frac{1}{n} \sum_{i=1}^n E \left[G_{ii}^2 \right] - \frac{2\lambda}{n + 2\lambda - 1}. \quad (73)$$

Noting finally that $E[G_{ii}^2] = s_i^2$ for $i \leq m$ (Eq. (31)) and $E[G_{ii}^2] = (n + 2\lambda + 1) / (n + 2\lambda - 1)$ for $i > m$ (from Eq. (25)), it is found that

$$\delta^2 = \frac{1}{n} \sum_{i=1}^m s_i^2 + \frac{n + 1 - (m/n)(n + 2\lambda + 1)}{n + 2\lambda - 1}. \quad (74)$$

5.2 Identification of the parameters λ , μ_{ii} , and τ_i

The above derivations have been carried out in terms of the parameters λ , μ_{ii} , and τ_i but these coefficients are not part of the original problem statement and thus they should, in principle, be evaluated in terms of the stated constraints, Eqs (5), (30), and (31). Note in this regard that the constraint of Eq. (5) has two aspects: the finiteness of \mathbf{v} and its specific value. The finiteness of \mathbf{v} guarantees the existence of the mean squared Frobenius norm of the inverse matrix \underline{A}^{-1} while the specific value of this coefficient providing an *overall* measure of the randomness of the matrices \underline{A} . In this light, \mathbf{v} and δ play a similar role and it is thus appropriate to replace Eq. (5) by a fixed value of δ , Eq. (74), see Soize (2001) for discussion. Proceeding in this manner would provide a direct expression for λ as

$$n + 2\lambda - 1 = 2q + 2 = \frac{n + 1 - 2(m/n)}{\left[\delta^2 - \frac{1}{n} \sum_{i=1}^m s_i^2 + \frac{m}{n} \right]}. \quad (75)$$

Note however that the available information about the randomness of the model lies in the standard deviations of some of its natural frequencies and thus it would be desirable to identify all model parameters from that information alone. Such an approach has already been proposed (Capiez-Lernout et al, 2005): the standard deviation of the first natural frequency was used to evaluate the parameter λ of an unconstrained model. A similar procedure is accordingly proposed here. Specifically, it will be assumed that the standard deviations of $m+1$ natural frequencies are known. Then, one of these standard deviation conditions will be used to yield the value of λ while the m others will be enforced through constraints of the form of Eq. (11).

To complete this strategy, it remains to determine which standard deviation should be singled out and used to compute λ . In this regard, it is valuable to analyze the behavior of the distribution of Eq. (39) as $q \rightarrow \infty$ as large values of this parameter occur for small to moderate system variability. In this limit, it is found that a positive value of τ_i is obtained only if

$$s_i^2 \leq 1 + 1/q. \quad (76)$$

In this regard, note that $1 + 1/q$ is the value of $E[G_{ii}^2]$ obtained from the unconstrained distribution of Eq. (25) in the same limit of $q \rightarrow \infty$. It is thus concluded that the specified value of s_i^2 must be less than what would be obtained by the unconstrained model with the same λ . This result is in fact not surprising since enforcing a higher value of $E[G_{ii}^2]$ would lead to an entropy larger than its maximum obtained for the unconstrained model.

This result also elucidates the determination of the parameter λ . Specifically, it will be evaluated from the standard deviation of the natural frequency which amongst the $m+1$ prescribed ones has the largest value of s_i^2 . Equality will be achieved in Eq. (76) for that value of i and the corresponding parameter τ_i will be zero.

The above observation that $E[G_{ii}^2] = 1 + 1/q$ as $q \rightarrow \infty$ independently of i for the unconstrained model implies first that the ratios of the standard deviation of a random natural frequency divided by the corresponding natural frequency of the mean model are all the same, i.e. $1/q$. Next, note that not imposing a second order constraint on the eigenvalues $m+1$ to n has led to the corresponding rows of the matrix \underline{H} being obtained from the unconstrained approach. Therefore, the natural frequencies $m+1$ to

n will then all have approximately the same ratio of standard deviation to value of mean model and equal to the largest s_i^2 . Equivalently, the natural frequencies $m+1$ to n would, in this manner, have their largest possible variances of value dictated by whichever eigenvalue exhibits the largest s_i^2 . While this condition will certainly lead to the largest value of the entropy, it would seem that the observed natural frequency of highest rank would provide the best basis for forecasting the variations of the unobserved, higher rank natural frequencies. Accordingly, it is proposed here to modify the modeling approach so that the parameters μ_{ii} and τ_i obtained for the row of \underline{H} associated with the observed natural frequency of highest rank be used also for the rows corresponding to unobserved, higher rank natural frequencies. Thus, the standard deviations of these random variables will be approximately (for finite value of q) constant as the rank increases.

6 NUMERICAL RESULTS AND VALIDATION

6.1 Evaluation of the Simulation Algorithm by Rejection

It is first desired to evaluate the performance, as measured by c_θ , of the simulation algorithm by rejection of the random variables G_{ii} , see section 4.3. The parameter c_θ , which represents the average number of pairs of uniform random numbers (U, \bar{U}) to be simulated before an acceptable sample G_{ii} is obtained, is determined by introducing Eq. (59) into Eq. (54). Note however that the resulting expression does depend on the normalization constant \bar{C}_i which could not be evaluated in closed form. Accordingly, the present validation will be carried first on limiting cases for which \bar{C}_i can be evaluated and then through a numerical study.

Previous applications of the nonparametric approach and the example below suggest that the parameter λ is generally much larger than 1 so that q , Eq. (53), is large. In such cases, the distribution of Eq. (39) can be closely approximated near its mode by a Gaussian distribution (see Devroye (1986) for a similar results for the Gamma distribution corresponding here to $\tau_i = 0$). To demonstrate this limiting behavior, note first that the mode y_0 tends to infinity as $q \rightarrow \infty$, see Eq. (52). Next, consider the behavior of $p_{G_{ii}}(y_0 + y)$ for fixed, finite values of y

as $y_0, q \rightarrow \infty$. With

$$(y_0 + y)^q = y_0^q \exp \left[q \ln \left(1 + \frac{y}{y_0} \right) \right] \approx y_0^q \exp \left[\frac{qy}{y_0} - \frac{qy^2}{2y_0^2} \right], \quad (77)$$

it is found that

$$p_{G_{ii}}(y_0 + y) \approx \hat{C}_i \exp \left[- \left(\tau_i + \frac{q}{2y_0^2} \right) y^2 \right] \quad (78)$$

where

$$\hat{C}_i = \bar{C}_i y_0^q \exp \left[-\mu_{ii} y_0 - \tau_i y_0^2 \right] \approx \sqrt{\frac{1}{\pi} \left(\tau_i + \frac{q}{2y_0^2} \right)}. \quad (79)$$

It is thus concluded from Eq. (67) that the distribution of Eq. (39) is, around its mode, approximately Gaussian. Next, combining Eqs (54), (59), and (79) provides the asymptotic limit of c_θ as $q \rightarrow \infty$ as

$$c_\theta = \sqrt{\frac{6}{\pi}} = 1.382 \quad (80)$$

for all values of μ_{ii} and τ_i . Note that the same asymptotic result was obtained from Eqs (60)-(64) thereby demonstrating the equivalency of Eqs (59) and (60)-(64) in the limit $q \rightarrow \infty$. It should finally be observed that Eq. (80) is valid independently of the mean and standard deviation of the approximate Gaussian distribution of Eq. (78) and thus is valid for all Gaussian distributions.

Although of a lesser practical interest, the other limiting situation of a small value of q was also investigated. Consider for example $q = 0$. In this case, the probability density function of Eq. (39) reduces to a truncated Gaussian distribution for which it is found that

$$\bar{C}_i = \sqrt{\frac{\tau_i}{\pi}} \exp \left(-\frac{\mu_{ii}^2}{4\tau_i} \right) \left[1 - \Phi \left(\frac{\mu_{ii}}{\sqrt{2}\tau_i} \right) \right]^{-1} \quad (81)$$

where $\Phi(\cdot)$ denotes the cumulative distribution function of a standard Gaussian random variable. The determination of the parameter c_θ is achieved separately for $\mu_{ii} \leq 0$ and $\mu_{ii} \geq 0$.

In the first case, the mode is $y_0 = -\mu_{ii} / 2\tau_i$ and the computation of c_θ proceeds from Eqs (59), or (60)-(64), and (54) and (81). For large values of $|\mu_{ii}| / \tau_i$, it is again found that the two values of c_θ obtained (one from Eq. (59) and the other from Eqs (60)-(64)) are very close to each other and to the value of Eq. (80). This result is not surprising as the truncation of the normal distribution plays a very small role for large values of the mode. As $|\mu_{ii}| / \tau_i$

decreases, the values of c_θ obtained by the two different algorithms differ increasingly from each other. For example, at $|\mu_{ii}| / \tau_i = 0$, the parameter c_θ corresponding to Eqs (60)-(64) converges to twice the value of Eq. (80), as expected since half of the produced Z values, those which are negative, are automatically rejected. On the contrary, the parameter c_θ obtained from Eq. (59) tends to infinity, due to the vanishing of the denominator of Eq. (59).

For positive values of μ_{ii} , there is no mode and Eq. (52) correctly yields the location of the maximum of the distribution as $y_0 = 0$. For $\mu_{ii} = 0$, the parameter c_θ is continuous and the values of $2\sqrt{6/\pi}$ and infinity are recovered for the algorithms based on Eqs (60)-(64) and (59), respectively. As μ_{ii} / τ_i increases, the two values of c_θ converge to each other and to 3 as $\mu_{ii} / \tau_i \rightarrow \infty$.

It is concluded from the above discussion that the algorithm based on Eqs (52), (54), (60)-(64) performs very well, typically requiring an average of 1.4 pairs of independent uniform random numbers per sample G_{ii} to be simulated. In unusual cases, the number of pairs of uniform variates may rise to 3. The algorithm based on Eq. (59) in place of Eqs (60)-(64) performs almost as well in the cases of greater practical interest but may lead to a large simulation effort when q is small (e.g. equal 0) and $|\mu_{ii}| / \tau_i$ is typically less than 3.

6.2 Numerical Estimation of the parameters λ , μ_{ii} , and τ_i

The initial conditions, μ_{ii}^0 and τ_i^0 , that were used to start these computations correspond to the large q Gaussian approximation of Eq. (39), i.e. Eqs (78) and (79). Accordingly, the constraints of Eqs (30) and (31) imply that

$$y_0^0 = \frac{1}{4\tau_i^0} \left(\sqrt{(\mu_{ii}^0)^2 + 8q\tau_i^0} - \mu_{ii}^0 \right) = 1 \quad (82)$$

and

$$s_i^2 - 1 = \frac{1}{2} \left(\tau_i^0 + \frac{q}{2(y_0^0)^2} \right)^{-1}. \quad (83)$$

The solution of Eqs (82) and (83) is readily found as

$$\tau_i^0 = \frac{1}{2} \left[\frac{1}{s_i^2 - 1} - q \right] \quad (84)$$

$$\mu_{ii}^0 = q - 2\tau_i^0 = 2q - \frac{1}{s_i^2 - 1}. \quad (85)$$

Next, it is necessary to evaluate μ_{ii} and τ_i from Eqs (30) and (31), or Eq. (30) and fixed second moments of the *exact* eigenvalues. In fact, neither of these constraints can be matched analytically given the complexity of Eq. (39) and thus a numerical approach must be undertaken. Either set of constraints (Eq. (30) and either Eq. (31) or the second moment of the exact eigenvalues) were considered and successfully matched using a polytope algorithm (IMSL routine DUMPOL).

At this point, note that the parameter τ_i^0 must be positive for the distribution of Eq. (39) to be integrable at infinity but this condition is automatically satisfied if the constraint of Eq. (76) is met.

6.3 Examples of Application

To exemplify the above developments, a $n=5$ degree of freedom dynamic system was considered exhibiting variations of its stiffness matrix. The above methodology was then applied with \underline{A} and \underline{B} being the stiffness and mass matrices, respectively. Further, the analysis of the system was carried out in the modal coordinates of the mean model with mass normalized modes. Thus, \underline{B} is the 5×5 unit matrix and \underline{A} is the diagonal matrix containing the squares of the natural frequencies which were first assumed to be 1, 3, 5, 7, and 9 rad/s (mean model 1). Finally, the damping matrix was constructed by assuming a damping ratio of 1% on all modes.

Three sets of computations were carried out, two with the unconstrained nonparametric approach and one with the present constrained formulation. The first unconstrained model was obtained by enforcing a standard deviation of the first natural frequency equal to 0.058, i.e. 5.8% of the corresponding value for the mean model. The second unconstrained model was similarly determined but with a lower standard deviation, i.e. 0.033. As stated earlier, it was observed that the ratios of the standard deviations of the natural frequencies of the random systems divided by their corresponding values for the mean models varied very little, from 0.00575 to 0.00584 for the first system and from 0.00331 to 0.00334 for the second one.

The constrained approach was applied next by imposing ratios of standard deviations to mean model value of 5.8% and 3.3% for the first and second natural frequency, respectively. Since the first natural frequency exhibits the largest relative variations, the first row of \underline{H} , i.e. the one most closely associated with the first frequency, was characterized by an unconstrained model, i.e. $\tau_1 = 0$. Given the lack of information on the third, fourth, and fifth

natural frequencies, the parameters μ_{ii} and τ_i were assumed to be the same for $i \geq 2$. Accordingly, there were only 4 model parameters to be determined: μ_{11} , λ , μ_{22} , and τ_2 . They were iteratively obtained by imposing the prescribed standard deviations on natural frequencies (the *exact* constraints, not the approximate ones of Eq. (11)) and Eq. (30) for $i=1$ and 2. In fact, Eq. (30) for $i=1$ is readily satisfied by selecting μ_{11} and λ to be related by Eq. (29) with $i=1$. Note further that the conditions of Eq. (30) for $i=3, 4$, and 5 are automatically satisfied when it is for $i=2$ as the model parameters μ_{ii} and τ_i are the same for $i \geq 2$. A sample of 10,000 random matrices was assumed for all computations.

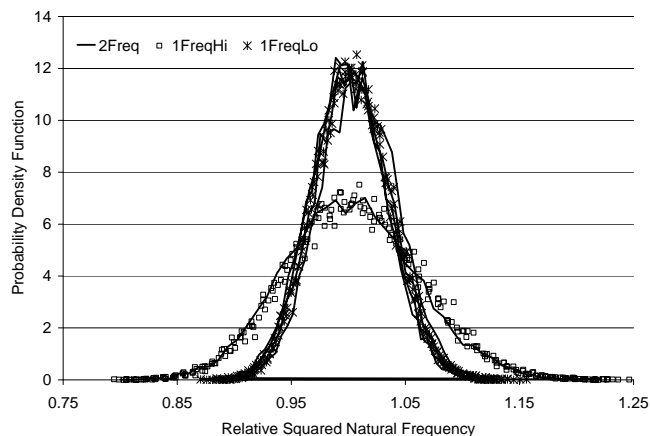


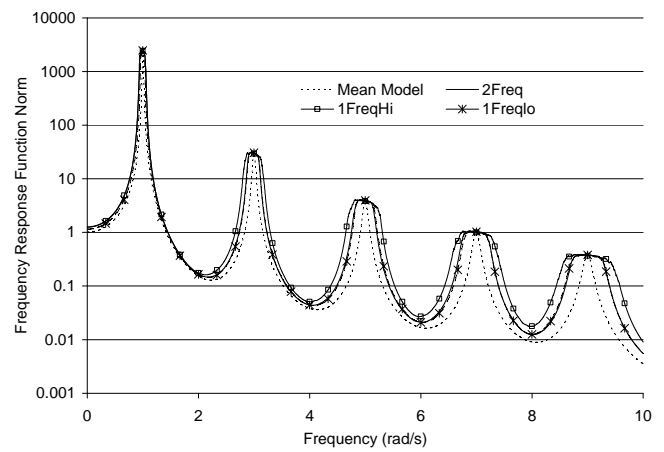
Figure 1. Probability density functions of the squared natural frequencies divided by their corresponding value for the mean model. Mean model 1. Constrained Model (“2Freq”), unconstrained models with higher (“1FreqHi”) and lower (“1FreqLo”) level of variations.

Shown in Fig. 1 are, for each of the three random systems, the distributions of the five natural frequencies divided by their corresponding value for the mean model. It is seen from this figure that the probability density function of the first natural frequency of the constrained system closely match its counterpart for the unconstrained system with the largest, i.e. 5.8%, level of variations. Similarly, the distributions of the 2nd-5th natural frequencies of the constrained system (divided by their corresponding value for the mean model) match those of the unconstrained model with the lowest, i.e. 3.3%, level of variations. It is thus concluded that the constrained model accurately matches the prescribed information on the variations of its natural frequencies.

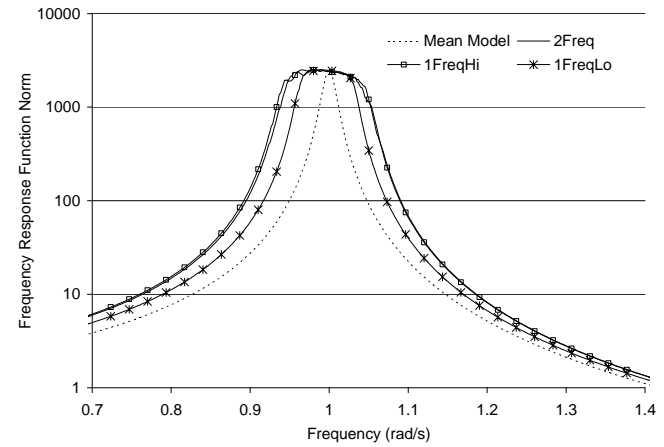
A comparison of the frequency response functions of the 3 random systems was carried out next. Shown in Fig. 2 are the 95th percentile of the Frobenius norm of the frequency response matrix

$\left[\underline{K} - \omega^2 \underline{M} + i\omega \underline{C} \right]^{-1}$ for the constrained model, for the two corresponding unconstrained models, and for the mean model. It is clearly seen from this figure (and additional results not shown here for brevity) that the frequency response functions of the constrained model and of the unconstrained one corresponding to the smaller (3.3%) variations in natural frequencies match very closely in the range $\omega \geq 1.5$ rad/s. For smaller frequencies however, e.g. see Fig. 2(b), the constrained model frequency response function closely match the unconstrained model with the largest variations natural in frequencies. These results are in fact in complete agreement with the distributions of natural frequencies presented in Fig. 1.

It was desired next to assess if the closeness of two natural frequencies of the mean model could affect the applicability of the approach and/or the validity of the above observations. To this end, the above computations were repeated exactly except for the second natural frequency of the mean model which was selected to be 1.06 rad/s., i.e. very close to the first one (mean model 2). The parameters of all three random systems were recomputed and the distributions of natural frequencies and norm of frequency response functions were again determined. Notwithstanding the closeness of the natural frequencies of the mean system, the constrained model accurately matched the prescribed variations of the natural frequencies, see Fig. 3. Note again that the distribution of the first natural frequency matches the one from the unconstrained model with higher level of variations, while the probability density function of the second one matches its counterpart for the other unconstrained model, as seen in Fig. 1. Finally, the norm of the frequency response function matrix, see Fig. 4, exhibits the properties already described in connection with Fig. 2. In regards to the behavior near 1Hz, note first that the two peaks of the mean model have merged when considering any of its statistics (mean, 5th, and 95th percentile) for either of the three random systems. Finally, note that the 95th percentile of the norm of the frequency response function of the constrained model matches closely its unconstrained counterpart with the larger level of variations left of the peak (as dictated by the first natural frequency) and the one with the lower level of variations right of the peak (as dictated by the second natural frequency).



(a)



(b)

Figure 2. (a) Frobenius norm of the frequency response function matrix, (b) same, zoomed. Mean model 1. Constrained Model (“2Freq”), unconstrained models with higher (“1FreqHi”) and lower (“1FreqLo”) level of variations.

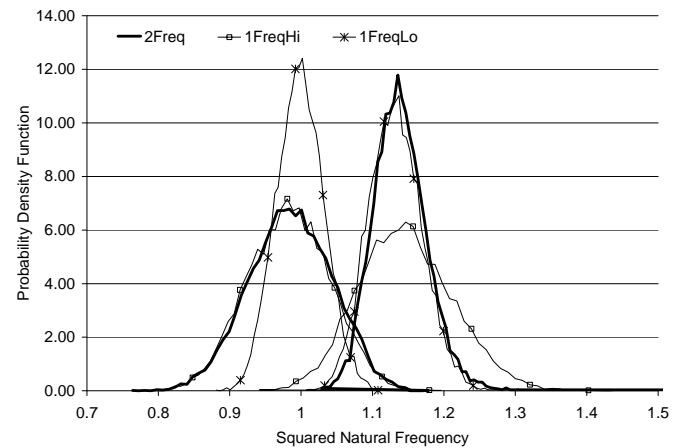
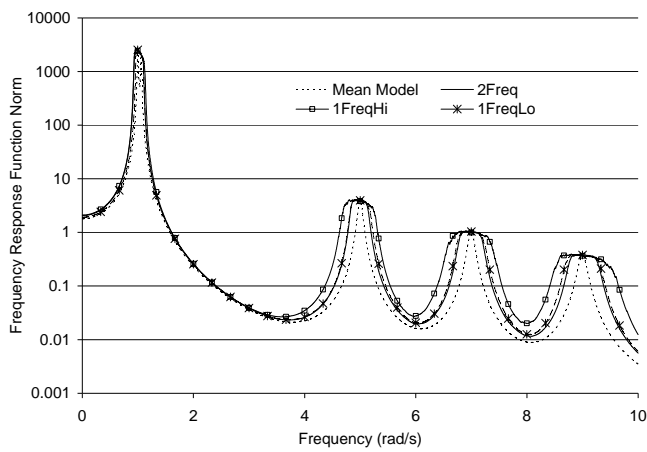
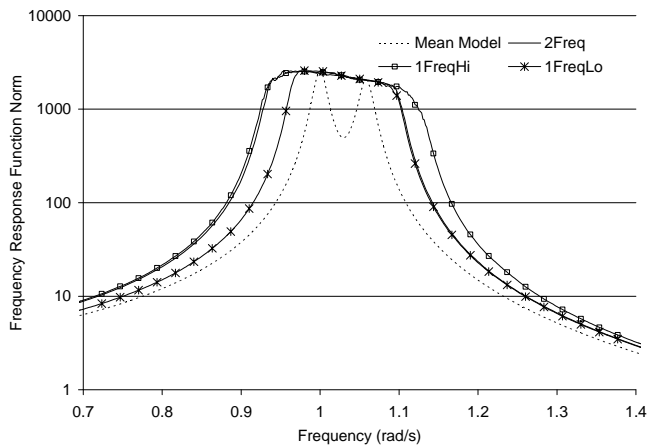


Figure 3. Probability density functions of the squared first and second natural frequencies. Mean model 2. Constrained Model (“2Freq”), unconstrained models with higher (“1FreqHi”) and lower (“1FreqLo”) level of variations.



(a)



(b)

Figure 4. (a) Frobenius norm of the frequency response function matrix, (b) same, zoomed. Mean model 2.Constrained Model (“2Freq”), unconstrained models with higher (“1FreqHi”) and lower (“1FreqLo”) level of variations.

7 SUMMARY

The focus of this paper was on the derivation and validation of a maximum entropy based probabilistic model of positive definite matrices, such as the stiffness and mass matrices of a dynamic system. Following prior work in this area, the maximum of the entropy is sought under the additional constraints that the mean matrix is known, that each realization $\underline{\underline{A}}$ is positive definite, and that the mean squared Frobenius norm of its inverse exists. A new set of constraints was however added in the present effort, Eq. (11), to permit the enforcement of prescribed standard deviations of an arbitrary number of natural frequencies of the random system. Accordingly, the present methodology provides a full probabilistic model of random systems from basic variability information. Further, this model is independent of a rotation/change of coordinate system, as desired. It should be noted that if only the standard deviation of a single natural frequency is prescribed, the present model reduces to the one previously developed.

Central to the inclusion of natural frequency variations in the model are the Rayleigh quotient type constraints of Eq. (11) and the representations of the mean and random matrices $\underline{\underline{A}}$ and $\underline{\underline{A}}$ in the specific form of Eqs (14), (15), and (17). These representations transfer the m frequency constraints of Eq. (11) into prescribed second order moments of the first m diagonal elements of the positive definite random matrix $\underline{\underline{G}}$, the mean value of which is identity. Enforcing the positive definiteness of $\underline{\underline{G}}$ is achieved by formulating directly the probabilistic model in terms of its Cholesky decomposition $\underline{\underline{H}}$.

In fact, the joint probability density function of its elements is given by Eq. (23) in which it is recognized that each row is stochastically independent of the others. Further, the first m rows, i.e. those associated with the natural frequency constraints, and the remaining ones are characterized by two different types of joint probability density functions. For rows $m+1$ to n , it is found, as in prior investigations, that all elements are independent of each other with the diagonal ones being distributed as square root of Gamma random variables. Further, the non-diagonal elements are found to follow a Gaussian distribution. The joint probability density function of the elements in the rows 1 to m is more complex because of the frequency constraints which induce a statistical coupling between these random variables. Nevertheless, the joint probability density function of these elements was obtained, Eq. (27), and discussed extensively. Further, a simple approach was given to simulate these elements which hinges on the generation of samples of the diagonal elements G_{ii} according to the distribution of Eq. (39). This task was achieved by an efficient, specially designed rejection algorithm, see Eqs (46) and (59)-(64).

The details of the application of this model were discussed and exemplified on the stiffness matrix of two 5 degree of freedom models, one exhibiting widely spread frequencies while the other has two very close natural frequencies. In both cases, the model accurately matched the prescribed information on the natural frequency variations and thus provides a sound and complete probabilistic description of the random stiffness matrix of these stochastic systems.

REFERENCES

- Capiez-Lernout, E., Soize, C., Lombard, J.-P., Dupont, C., and Seinturier, E. 2005. Blade manufacturing tolerances definition for a mistuned industrial bladed disk, *Journal of Engineering for Gas Turbines and Power*, Vol. 127, No. 3, pp. 621-628.

- Devroye, L. 1986. Non-Uniform Random Variate Generation, Springer Verlag, New York.
- Ghanem, R., and Spanos, P.D. 1991. Stochastic Finite elements: A spectral approach, Spring-Verlag, New York.
- Kleiber, M., Tran, D.H., Hien, T.D. 1992. The Stochastic Finite element Method, John Wiley and Sons, new York, 1992.
- Rivas-Guerra, A.J., and Mignolet, M.P. 2004. Local/Global Effects of Mistuning on the Forced Response of Bladed Disks, *Journal of Engineering for Gas Turbines and Power*, Vol. 126, No. 1, pp. 131-141.
- Schueller, G.I. (Ed.). 1997. A State-of-the-Art Report on Computational Stochastic Mechanics, *Probabilistic Engineering Mechanics*, Vol. 12, No. 4, pp. 197-321.
- Soize, C. 2000. A nonparametric model of random uncertainties on reduced matrix model in structural dynamics, *Probabilistic Engineering Mechanics*, Vol. 15, No. 3, pp. 277-294.
- Soize, C. 2001. Maximum entropy approach for modeling random uncertainties in transient elastodynamics, *Journal of the Acoustical Society of America*, Vol. 109, No. 5, pp. 1979-1996.
- Soize, C. 2005. Probabilistic models for computational stochastic mechanics and applications, *Proceedings of the 9th International Conference on Structural Safety and Reliability ICOSSAR'05, Rome, Italy, 19-23 June 2005*, Augusti, G., Schueller, G.I., and Ciampoli, M., (Eds). Millpress, Rotterdam, Netherlands, pp.23-42.