MODEL SELECTION IN THE SPACE OF GAUSSIAN MODELS INVARIANT BY SYMMETRY

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We consider multivariate centered Gaussian models for the random variable \( Z = (Z_1, \ldots, Z_p) \), invariant under the action of a subgroup of the group of permutations on \( \{1, \ldots, p\} \). Using the representation theory of the symmetric group on the field of reals, we derive the distribution of the maximum likelihood estimate of the covariance parameter \( \Sigma \) and also the analytic expression of the normalizing constant of the Diaconis-Ylvisaker conjugate prior for the precision parameter \( K = \Sigma^{-1} \). We can thus perform Bayesian model selection in the class of complete Gaussian models invariant by the action of a subgroup of the symmetric group, which we could also call complete RCOP models. We illustrate our results with a toy example of dimension 4 and several examples for selection within cyclic groups, including a high dimensional example with \( p = 100 \).

1. Introduction.

1.1. Motivations and applications. Let \( V = \{1, \ldots, p\} \) be a finite index set and let \( Z = (Z_1, \ldots, Z_p) \) be a multivariate random variable following a centered Gaussian model \( N_p(0, \Sigma) \). Let \( \mathfrak{S}_p \) denote the symmetric group on \( V \), that is, the group of all permutations on \( \{1, \ldots, p\} \) and let \( \Gamma \) be a subgroup of \( \mathfrak{S}_p \). A centered Gaussian model is said to be invariant under the action of \( \Gamma \) if for all \( g \in \Gamma \), \( g \cdot \Sigma \cdot g^\top = \Sigma \) (here we identify a permutation \( g \) with its permutation matrix).

Given \( n \) data points \( Z^{(1)}, \ldots, Z^{(n)} \) from a Gaussian distribution, our aim in this paper is to do Bayesian model selection within the class of models invariant by symmetry, that is, invariant under the action of some subgroup \( \Gamma \) of \( \mathfrak{S}_p \) on \( V \). Given the data, our aim is therefore to identify the subgroup \( \Gamma \subset \mathfrak{S}_p \) such that the model invariant under \( \Gamma \) has the highest posterior probability. We achieve this goal by constructing a Markov chain on the space of models and using the Metropolis-Hastings algorithm.

There are many alternative ways of doing model search in modern statistics on big data sets, both frequentist and Bayesian. Bayesian model selection methods (cf. (Ghosal and van der Vaart, 2017, Ch.10)) are widely used in practice, thanks to the possibility of using a prior knowledge on the model and to their rigorous mathematical bases. Moreover, in

*Supported by JSPS KAKENHI Grant Number 16K05174 and JST PRESTO.
†Supported by grant 2016/21/B/ST1/00005 of the National Science Center, Poland.
‡Supported by an NSERC Discovery Grant.


Keywords and phrases: colored graph, conjugate prior, covariance selection, invariance, permutation symmetry.
Bayesian approach the Metropolis-Hastings algorithm is naturally applicable and generally accepted.

Our work can be viewed as a special case of colored graphical Gaussian models (the underlying graph is complete so we do not impose conditional independence structure), that is, statistical graphical models with additional symmetries (equality constraints on the precision or correlation matrix). Such models were introduced into modern exploratory analysis of data in the seminal paper Højsgaard and Lauritzen (2008), as a powerful tool of dimension reduction in unsupervised learning, cf. (Maathuis et al., 2018, Chapter 9.8). A preponderant role is given in Højsgaard and Lauritzen (2008) to the RCOP models studied in our paper, thanks to their most tractable structure and interpretability through symmetries among the variables. One of motivations of this paper is to address the task stated in (Højsgaard and Lauritzen, 2008, p. 1025): For the models to become widely applicable, it is mandatory to develop algorithms for model identification which are robust, reliable and transparent.

For high-dimensional data, Gaussian models which have symmetries and are graphical allow statisticians to reduce the dimension of a model. In genetics, such models can be used to identify genes having the same function or groups of genes having similar interactions. Below we mention some of the studies in which our model could find potential application.

In Højsgaard and Lauritzen (2008), gene expression signatures for p53 mutation status in 58 breast cancer samples consisting of 150 genes were investigated and interpreted. We apply our algorithm to this data in the Supplementary material, see Section 6 in Graczyk et al.. We claim that our model selection procedure can be used as an exploratory tool. Assuming that the variables are all on some common scale, the proposed algorithm can be run to look for potential hidden symmetries between the variables.

It is worth to underline that one of the characteristics of our model is the lack of scale invariance. We point out below that there are many examples where our model can still be applied. E.g. the data from gene expression are on the same scale in the sense that they are results of experiments of the same type, measured in the same gauges. Similar situation appears generally for omic data sets in proteomics and metabolomics. For more details see e.g. the monograph Frommlet, Bogdan and Ramsey (2016). In (Sobczyk et al., 2020, Section 6.2 TCGA Breast Cancer Data), genetic information in tumoral tissues DNA that are involved in gene expression are measured from messenger sequencing by the RNASeq method and they are all on the same scale, as they are the numbers of transcripts in a sample. In clinical epidemiology and medicine, one often uses scales combined into scores to classify outcome, see e.g. Toyoda et al. (2021); Missio et al. (2019). Range of values of such scores are often similar, even though not formally tested statistically to be so. In the paper Descatha et al. (2007) the normalization or non-normalization of data did not influence their statistical interpretation.

Moreover, we argue that it is natural to expect certain symmetries in the data from gene expression. Namely, expression of a given gene is triggered by binding the transcription factors to the gene transcription factor binding sites. The transcription factors are the proteins produced by other genes, say regulatory genes. In the gene network there are often many genes triggered by the same regulatory genes and it makes sense to assume that their relative expressions depend on the abundance of proteins of the regulatory genes (i.e. gene expressions) in a similar way.

In Gao and Massam (2015), 12 625 neutrophil gene expressions were monitored with imposed symmetry constraints to the graphical modeling. The paper Li, Gao and Massam (2020) contains a study of the structure of colored graphs applied to a flow cytometry data set on signaling networks of human immune system cells, which consists of 7 466 measurements on 11 phosphorylated proteins.
A very recent application of graphical models with symmetries to fMRI real data on brain networks is proposed in Ranciati, Roverato and Luati (2021). An impressive number of recent applications of graphical models to real data analysis is listed in the recent monograph (Maathuis et al., 2018, Chapters 19,20,21) and includes genetics, genomics, molecular systems biology and forensic analysis, cf. also the books Roverato (2017) for medical and Li (2009) for image data applications.

Finally, let us mention that colored graphical models provide interesting examples of exponential algebraic varieties and algebraic exponential families, e.g. Toeplitz matrices Michałek et al. (2016); see also Davies and Marigliano (2021). The recent algebro-geometric approach to graphical models and Gaussian Bayesian networks is being developed intensely (Maathuis et al., 2018, Chapter 3).

1.2. Contribution of the paper and relations to previous work. In this subsection we carefully describe and position this paper in the context of previous research.

Theory of invariant normal models (with the so-called lattice conditional independencies Andersson and Madsen (1998); Madsen (2000), which are not considered in the present paper) was developed by the Danish school. History regarding this subject is nicely presented in Andersson and Madsen (1998), where the reader can also find references to earlier works dealing with particular symmetry models such as, for example, the circular symmetry model of Olkin and Press (1969) that we will consider further (Section 5). Among others, Andersson, Brøns, Jensen, Madsen and Perlman developed a fairly complete theory of MLE $\hat{\Sigma}$ of the covariance matrix in invariant normal models, however, the problems considered in our paper are very different.

These works were concentrating on the derivation of statistical properties of the maximum likelihood estimate of $\Sigma$ and on testing the hypothesis that models were of a particular type. In particular, to the best of our knowledge the Danish school never considered any model search in the context of invariant normal models. When the model space is very big (and this is the usual case of our framework), then it is impossible to perform simultaneous tests for all possible models. Despite the computation problems, there is also even bigger issue due to multiple comparisons problem.

Just like the classical papers mentioned above, the fundamental algebraic tool we use in this work is the irreducible decomposition theorem for the matrix representation of the group $\Gamma$, which in turn means that, through an adequate change of basis, any matrix $X$ in $\mathbb{Z}_\Gamma$, the space of symmetric matrices invariant under the subgroup $\Gamma$ of $\mathfrak{S}_p$, can be written in a block diagonal form. The following result is a reformulation of an observation made in (Andersson, 1975, 4.6-4.8).

**Theorem 1.** Fix a permutation subgroup $\Gamma \subset \mathfrak{S}_p$. Then, there exist constants $L \in \mathbb{N}$, $(k_i, d_i, r_i)_{i=1}^L$ and orthogonal matrix $U_\Gamma$ such that if $X \in \mathbb{Z}_\Gamma$, i.e. $X \in \text{Sym}(p; \mathbb{R})$ and $X_{ij} = X_{\sigma(i)\sigma(j)}$ $(\sigma \in \Gamma$, $i,j \in \{1, \ldots, p\})$, then $X = U_\Gamma \cdot \begin{pmatrix} M_{K_1}(x_1) \otimes I_{k_1/d_1} & M_{K_2}(x_2) \otimes I_{k_2/d_2} & \cdots \\ M_{K_L}(x_L) \otimes I_{k_L/d_L} \end{pmatrix} \cdot U_\Gamma^T$, where $M_{K_i}(x_i)$ is a real matrix representation of an $r_i \times r_i$ Hermitian matrix $x_i$ with entries in $K_i = \mathbb{R}, \mathbb{C}$ or $\mathbb{H}$, $i = 1, \ldots, L$, and $A \otimes B$ denotes the Kronecker product of matrices $A$ and $B$. 
Elements of \((k_i, d_i, r_i)_{i=1}^L\) are integer constants called structure constants that we will define later. At this point we note that \(k_i/d_i\) are also integers and \(d_i = \text{dim}_\mathbb{R} \mathbb{K}_i \in \{1, 2, 4\}\). The mappings \(M_{K_i} : \text{Herm}(r_i; \mathbb{K}_i) \to \text{Sym}(d_i r_i; \mathbb{R})\) are defined in Section 2.2. As was already observed in Jensen (1988), the space \(Z_\Gamma\) equipped with a Jordan product and trace inner product forms a Euclidean Jordan algebra. Thus, (1) can be understood as a decomposition of \(Z_\Gamma\) into Euclidean simple Jordan algebras. Theorem 1 is the existence result and actual computation of structure constants and the orthogonal matrix \(U_\Gamma\) is in general a hard technical task. A complete proof of Theorem 1 can be found in the Supplementary material Graczyk et al.. We tried to ensure that our arguments are concrete and should be easier to understand for the reader who is not familiar with representation theory.

The main novel results of the paper are

\(\mathbf{a}\) new Bayesian model selection procedure within Gaussian models invariant by a permutation subgroup, Section 4.1,

\(\mathbf{b}\) explicit formulas for Gamma integrals, normalizing constants of densities of Diaconis-Ylvisaker conjugate prior for \(K\) and of the MLE of \(\Sigma\) on \(\mathcal{P}_\Gamma\), Bayes factors, which are necessary for performing \(\mathbf{a}\), Theorems 8 and 9 in Section 3,

\(\mathbf{c}\) efficient algorithm for finding a decomposition (1) when the subgroup \(\Gamma\) is cyclic, Theorems 5 and 6 in Section 2.4,

\(\mathbf{d}\) simulations that visualize the performance of the method in low and high dimensional examples, Section 4.2, Section 5 and Section 4 of the Supplement Graczyk et al..

\textbf{Ad (a)}. We are aware of three papers which concern model selection in the space of colored graphical model, namely Gehrmann (2011); Massam, Li and Gao (2018); Li, Gao and Massam (2020).

In Gehrmann (2011) the author used the lattice structure of the colored graphical model classes and applied Edwards–Havránek model selection procedure to \(p = 4\) and \(p = 5\) examples, admitted that applying this method to high-dimensional problems requires additional work.

Both papers Massam, Li and Gao (2018) and Li, Gao and Massam (2020) used Bayesian methods and allow for model selection in the space of RCON models (which is a superclass of RCOP models introduced in Højsgaard and Lauritzen (2008)) and for arbitrary graphs describing conditional independencies in a vector. Such generality comes at a certain cost: as the authors were not able to compute normalizing constants for such general models, they had to approximate these constants or bypass the problem (which comes with a significant increase in computational complexity); we quote a few lines from these articles that describe the situation well.

- Massam, Li and Gao (2018): \textit{However, just as sampling schemes for the \(G\)-Wishart distribution are not recommended for computation of (normalizing constant) \(I_G(\delta, D)\) and model selection in higher dimensions, our sampling scheme is not recommended for computing (normalizing constant) \(I_G(\delta, D)\) in high dimensions.}

- Li, Gao and Massam (2020): \textit{The model \(G^*\) with an additional edge is then compared to the current model \(G\) using the Bayes factor (...) which itself is computed with the help of the double reversible jump MCMC algorithm. (...) We thus avoid computing these quantities which are the usual computational stumbling blocks in graphical Gaussian model selection.}

Our approach to the Bayesian model selection is much simpler as we were able to compute normalizing constants of Diaconis-Ylvisaker conjugate priors for \(K\).

\textbf{Ad (b)}. We note that a general form of a density of the MLE under our assumptions was already written in Andersson (1975) and in more explicit form in Andersson and Madsen (1998). However, an explicit expression for the normalizing constant of density of \(\hat{\Sigma}\) or
Diaconis-Ylvisaker conjugate prior was not the object of interest of the Danish school and it is crucial for the Bayes paradigm and the Bayesian model selection.

Still, there are certain results in their numerous works that can be compared with our formulas. In particular, (Andersson and Madsen, 1998, Eq. (A.4)) gives a formula for $\mathbb{E}[\text{Det}(\hat{\Sigma})^{\alpha}]$ which is consistent with our results. Indeed, after substitution of $(d_i, n_i, p_i)$ for $(d_i, n_i, r_i)$, the right hand side of their formula coincides with $2^{mp} n^{-mp} \Gamma_p (\alpha + n/2) / \Gamma_p (n/2)$ in our notation (see Theorem 8). Further, in (Andersson, Brøns and Jensen, 1983, Section 8) explicit formula for normalizing constants of the density of eigenvalues of $\hat{\Sigma}$ is given. However, as distribution of eigenvalues of a random matrix does not determine the distribution of this matrix, our formulas do not follow from these results.

In some very special cases, normalizing constants for Diaconis-Ylvisaker conjugate prior are given in Massam, Li and Gao (2018).

**Ad (c).** In order to compute normalizing constants in our model, one needs to know explicit decomposition (1), that is, the structure constants and the orthogonal matrix $U_\Gamma$. The same issue can be seen in (Jensen, 1988, Theorem 1), which is the existence result (like our Theorem 1) and does not give the answer how should one proceed to find such decomposition. In order for this theory to be applied, we proved that when $\Gamma$ is a cyclic subgroup, then we can efficiently find explicit decomposition (1) for arbitrary $p$. This practical aspect of our work has not been addressed before. To our knowledge, our paper is the first one to identify a non-trivial class of subgroups for which all objects can be calculated explicitly.

For a moment, let us consider the more general situation of Gaussian graphical models with conditional independence structure encoded by a non complete graph $G$. Then one can introduce symmetry restrictions (RCOP) by requiring that the precision matrix $K$ is invariant under some subgroup $\Gamma$ of $\mathfrak{S}_p$. However, when $G$ is not complete, not all subgroups are suited to the problem. In such cases, one has to require that $\Gamma$ belongs to the automorphism group $\text{Aut}(G)$ of $G$. If a graph $G$ is sparse, then $\text{Aut}(G)$ may be very small and it is natural to expect that the vast majority of subgroups of $\text{Aut}(G)$ are actually cyclic. Moreover, finding the structure constants for a general group is much more expensive and in some situations it may not be worth to consider the problem in its full generality. We consider our work as a first step towards the rigorous analytical treatment of Bayesian model selection in the space of graphical Gaussian models invariant under the action of $\Gamma \subset \mathfrak{S}_p$ when conditional independencies are allowed.

Moreover, we offer here a new heuristic approach to colored graphical models using our “full graph” approach. It was already observed in Højsgaard and Lauritzen (2008) that the color pattern of the covariance matrix and the precision matrix are the same (i.e. they belong to the space $Z_\Gamma$). The same applies to the off-diagonal elements of the partial correlation matrix. Our procedure allows one to find the color pattern of the covariance matrix. Since our model does not suppose any preliminary conditional independence structure, the corresponding graph is complete and there are no zeros in the partial correlation matrix. However, if the true graph is not complete, it is natural to expect from the model that similar entries of the partial correlation matrix (in particular those which are close to 0) are colored in the same way. Thus, to recover the true graph we may threshold the values of the partial correlation matrix. More precisely, we choose a threshold $\alpha > 0$ and we construct a colored graph $G$ by maintaining the color pattern previously found and requiring that for $i \neq j$,

$$i \sim j \quad \text{if and only if} \quad \frac{|k_{ij}|}{\sqrt{k_{ii}k_{jj}}} > \alpha,$$

where $K = (k_{ij})_{i,j}$ is the precision matrix. Resulting graph $G$ is in general not complete and the corresponding space of admissible covariance matrices is still invariant under the action
of the subgroup found by our procedure; thus we obtain a RCOP model. We applied this approach to a real data example in Section 4 of the Supplementary material.

There are also several recent papers which use a version of Theorem 1. The subject of Soloveychik, Trushin and Wiesel (2016) is estimation of complex covariance matrices in complex random vectors in non-Gaussian models invariant under the action of a fixed permutation subgroup, see also De Maio et al. (2016). We remark that the argument of Soloveychik, Trushin and Wiesel (2016) is based on representation theory over complex number fields, and as was noticed by them, the fundamental structure theorem is much simpler than Theorem 1 because of the difference between the representation theory over \( \mathbb{C} \) and \( \mathbb{R} \). In Shah and Chandrasekaran (2012) the authors consider the real case and sub-Gaussian model for which they establish rates of convergence of an estimator of \( \Sigma \), empirical covariance matrix regularized by the action of a known permutation subgroup.

1.3. Outline of the paper. Let us consider the following example, which shows how Theorem 1 works.

**Example 1.** For \( p = 3 \) and \( \Gamma = \mathfrak{S}_3 \), the space of symmetric matrices \( X \) invariant under \( \Gamma \), that is, such that \( X_{ij} = X_{\sigma(i)\sigma(j)} \) for all \( \sigma \in \Gamma \), is

\[
\mathcal{Z}_\Gamma = \left\{ \begin{pmatrix} a & b \\ b & a \end{pmatrix} ; a, b \in \mathbb{R} \right\}.
\]

The decomposition (1) yields \( U_\Gamma := (v_1, v_2, v_3) \in O(3) \) with

\[
v_1 := \begin{pmatrix} 1/\sqrt{3} \\ 1/\sqrt{3} \\ 1/\sqrt{3} \end{pmatrix}, \quad v_2 := \begin{pmatrix} \sqrt{2}/3 \\ -1/\sqrt{6} \\ -1/\sqrt{6} \end{pmatrix}, \quad v_3 := \begin{pmatrix} 0 \\ 1/\sqrt{2} \\ -1/\sqrt{2} \end{pmatrix},
\]

and

\[
\begin{pmatrix} a & b \\ b & a \end{pmatrix} = U_\Gamma \cdot \begin{pmatrix} a + 2b & a - b \\ a - b & a - b \end{pmatrix} \cdot U_\Gamma^T.
\]

Here \( L = 2 \), \( k_1/d_1 = 1 \), \( k_2/d_2 = 2 \), \( \mathbb{K}_1 = \mathbb{K}_2 = \mathbb{R} \), \( d_1 = d_2 = 1 \).

We see immediately in the example above that, following the decomposition (1), the trace \( \text{Tr} \[X\] = a + 2b + 2(a - b) \) and the determinant \( \text{Det} \[X\] = (a + 2b)(a - b)^2 \) can be readily obtained. Similarly, using (1) allows us to easily obtain \( \text{Det} \[X\] \) and \( \text{Tr} \[X\] \) in general.

In Section 3, we will see that having the explicit formulas for \( \text{Det} \[X\] \) and \( \text{Tr} \[X\] \), in turn, allows us to derive the analytic expression of the Gamma function on \( \mathcal{P}_\Gamma = \mathcal{Z}_\Gamma \cap \text{Sym}^+(p; \mathbb{R}) \), defined as

\[
\Gamma_{\mathcal{P}_\Gamma}(\lambda) := \int_{\mathcal{P}_\Gamma} \text{Det} \[X\]^\lambda e^{-\text{Tr}[X]} \varphi_{\Gamma}(X) \, dX,
\]

where \( \varphi_{\Gamma}(X) \, dX \) is the invariant measure on \( \mathcal{P}_\Gamma \) (see Definition 10 and Proposition 7) and \( dX \) denotes the Euclidean measure on the space \( \mathcal{Z}_\Gamma \) with the trace inner product.

With our results, we can derive the analytic expression of the normalizing constant \( I_{\Gamma}(\delta, D) \) of the Diaconis-Ylvisaker conjugate prior on \( K = \Sigma^{-1} \) with density, with respect to the Euclidean measure on \( \mathcal{Z}_\Gamma \), equal to

\[
f(K; \delta, D) = \frac{1}{I_{\Gamma}(\delta, D)} \text{Det} \[K\]^{(\delta - 2)/2} e^{-\frac{1}{2} \text{Tr}[K \cdot D]} \mathbf{1}_{\mathcal{P}_\Gamma}(K) \]
for appropriate values of the scalar hyper-parameter $\delta$ and the matrix hyper-parameter $D \in \mathcal{P}_T$. By analogy with the $G$-Wishart distribution, defined in the context of the graphical Gaussian models, Markov with respect to an undirected graph $G$ on the cone $\mathcal{P}_G$ of positive definite matrices with zero entry $(i, j)$ whenever there is no edge between the vertices $i$ and $j$ in $G$, (see Maathuis et al. (2018)), we can call the distribution with density $f(K; \delta, D)$, the RCOP-Wishart (RCOP is the name coined in Højsgaard and Lauritzen (2008) for graphical Gaussian models with restrictions generated by permutation symmetry). It is important to note here that if $\Sigma$ is in $\mathcal{P}_T$, so is $K = \Sigma^{-1}$ so that $K$ can also be decomposed according to (1). Equipped with all these results, we compute the Bayes factors comparing models pairwise and perform model selection. We will indicate in Section 4 how to travel through the space of subgroups of the symmetric group.

In Section 3, we also derive the distribution of the maximum likelihood estimate (henceforth abbreviated MLE) of $\Sigma$ and show that for $n \geq \max_{i=1,\ldots,L} \{ r_i d_i / k_i \}$ it has a density equal to

$$\frac{\text{Det}(X)^{n/2} e^{-\frac{1}{2} \text{Tr}[X^{-1} \Sigma^{-1}]} \varphi_{\Gamma}(X) \mathbb{1}_{\mathcal{P}_T}(X)}{\text{Det}(2\Sigma)^{n/2} \Gamma_{\mathcal{P}_T}(\frac{\sigma}{2})}.$$

Clearly, the key to computing the Gamma integral on $\mathcal{P}_T$, the normalizing constant $I_{\Gamma}(\delta, D)$ or the density of the MLE of $\Sigma$ is, for each $\Gamma \subset S_p$, to obtain the block diagonal matrix with diagonal block entries $M_{\Gamma}(x_i) \otimes I_{k_i/d_i}, i = 1, \ldots, L$, in the decomposition (1). In principle, we have to derive the invariant measure $\varphi_{\Gamma}$ and find the structure constants $(k_i, d_i, r_i)_{i=1}^L$. This goal can be achieved by constructing an orthogonal matrix $U_{\Gamma}$ and using (1). However, doing so for every $\Gamma$ visited during the model selection process is computationally heavy.

We will show that for small to moderate dimensions, we can obtain the structure constants as well as the expression of $\text{Det}(X)$ and $\varphi_{\Gamma}(X)$ without having to compute $U_{\Gamma}$. Indeed, as indicated in Lemma 4, for any $X \in \mathcal{P}_T$, $\text{Det}(X)$ admits a unique irreducible factorization of the form

$$\text{Det}(X) = \prod_{i=1}^L \text{Det}(M_{\Gamma_i}(x_i))^{k_i/d_i} = \prod_{j=1}^L f_j(X)^{a_j}, \quad (X \in \mathcal{Z}_T),$$

where each $a_j$ is a positive integer, each $f_j(X)$ is an irreducible polynomial of $X \in \mathcal{Z}_T$, and $f_i \neq f_j$ if $i \neq j$. The constants $k_i, d_i, r_i$ are obtained by identification of the two expressions of $\text{Det}(X)$ in (2). Factorization of a homogeneous polynomial $\text{Det}(X)$ can be performed using standard software such as either Mathematica or Python.

Due to computational complexity, for bigger dimensions, it is difficult to obtain the irreducible factorization of $\text{Det}(X)$. For special cases such as the case where the subgroup $\Gamma$ is a cyclic group, we give (Section 2.4) a simple construction of the matrix $U_{\Gamma}$ and thus, for any dimension $p$, we can do model selection in the space of models invariant under the action of a cyclic group. We argue that restriction to cyclic groups is not as limiting as it may look. The formula for the number of different colorings $c_p = \# \{ \mathcal{Z}_T : \Gamma \subset S_p \}$ for given $p$ is unknown. Obviously, it is bounded from above by the number of all subgroups of $S_p$, because different subgroups may produce the same coloring (e.g. in Example 1 we have $\mathcal{Z}_{S_5} = \mathcal{Z}_{\{1,2,3\}}$). On the other hand, it is known (see Lemma 15) that $c_p$ is bounded from below by the number of distinct cyclic subgroups, which grows rapidly with $p$ (see OEIS1 sequence A051625). In particular, for $p = 18^2$, we have $c_p \in (7.1 \cdot 10^{14}, 7.6 \cdot 10^{18})$, see also Table 1. The lower bound

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2 The number of subgroups of $S_p$ is unknown for $p > 18$, see Holt (2010) and OEIS sequence A005432.
for $c_p$ indicates that the colorings obtained from cyclic subgroups form a rich subfamily of all possible colorings.

The procedure to do model selection will be described in Section 4 and we will illustrate this procedure with Frets’ data (see Frets (1921)) and several examples for selection within cyclic groups, including a high dimensional example with $p = 100$ (Section 5) and a real data example Miller et al. (2005) with $p = 150$ in Section 4 of the Supplement.

2. Preliminaries and structure constants. In this section we present methods to calculate the structure constants of a decomposition given in Theorem 1. Additions to this section can be found in Section 3 of the Supplementary material Graczyk et al..

2.1. Notation. Let $\text{Mat}(n_m; \mathbb{R})$, $\text{Sym}(n_\mathbb{R})$ denote the linear spaces of real $n \times m$ matrices and symmetric real $n \times n$ matrices, respectively. Let $\text{Sym}^+(n_\mathbb{R})$ be the cone of symmetric positive definite real $n \times n$ matrices. $A^\top$ denotes the transpose of a matrix $A$. Det and $\text{Tr}$ denote the usual determinant and trace in $\text{Mat}(n, n; \mathbb{R})$.

For $A \in \text{Mat}(m, n; \mathbb{R})$ and $B \in \text{Mat}(m', n'; \mathbb{R})$, we denote by $A \oplus B$ the matrix $\begin{pmatrix} A & O \\ O & B \end{pmatrix} \in \text{Mat}(m + m', n + n'; \mathbb{R})$, and by $A \otimes B$ the Kronecker product of $A$ and $B$. For a positive integer $r$, we write $B^{\otimes r}$ for $I_r \otimes B \in \text{Mat}(rm', rnm'; \mathbb{R})$.

Let $p$ denote the fixed number of vertices of a graph and let $\mathfrak{S}_p$ denote the symmetric group. We write permutations in cycle notation, meaning that $(i_1, i_2, \ldots, i_n)$ maps $i_j$ to $i_{j+1}$ for $j = 1, \ldots, r - 1$ and $i_n$ to $i_1$. By $\langle \sigma_1, \ldots, \sigma_k \rangle$ we denote the group generated by permutations $\sigma_1, \ldots, \sigma_k$. The composition (product) of permutations $\sigma, \sigma' \in \mathfrak{S}_p$ will be denoted by $\sigma \circ \sigma'$.

DEFINITION 2. For a subgroup $\Gamma \subset \mathfrak{S}_p$, we define the space of symmetric matrices invariant under $\Gamma$, or the vector space of colored matrices,

$$Z_{\Gamma} := \{ x \in \text{Sym}(p; \mathbb{R}) ; x_{ij} = x_{\sigma(i)\sigma(j)} \text{ for all } \sigma \in \Gamma \},$$

and the cone of positive definite matrices valued in $Z_{\Gamma}$,

$$P_{\Gamma} := Z_{\Gamma} \cap \text{Sym}^+(p; \mathbb{R}).$$

We note that the same colored space and cone can be generated by two different subgroups: in Example 1, the subgroup $\Gamma' = \langle (1, 2, 3) \rangle$ generated by the permutation $\sigma = (1, 2, 3)$ is such that $\Gamma' \neq \Gamma$ but $Z_{\Gamma'} = Z_{\Gamma}$. Let us define

$$\Gamma^* = \{ \sigma^* \in \mathfrak{S}_p ; x_{ij} = x_{\sigma^*(i)\sigma^*(j)} \text{ for all } x \in Z_{\Gamma} \}.$$ 

Clearly, $\Gamma$ is a subgroup of $\Gamma^*$ and $\Gamma^*$ is the unique largest subgroup of $\mathfrak{S}_p$ such that $Z_{\Gamma^*} = Z_{\Gamma}$ or, equivalently, such that the $\Gamma^*$- and $\Gamma$- orbits in $\{ \{ v_1, v_2 \} ; v_i \in V, i = 1, 2 \}$ are the same. The group $\Gamma^*$ is called the 2*-closure of $\Gamma$. The group $\Gamma$ is said to be 2*-closed if $\Gamma = \Gamma^*$. Subgroups which are 2*-closed are in bijection with the set of colored spaces. These concepts have been investigated in Wielandt (1969); Siemons (1982) along with a generalization to regular colorings in Siemons (1983). The combinatorics of 2*-closed subgroups is very complicated and little is known in general. (Graham, Grötschel and Lovász, 1995, p. 1502). In particular, the number of such subgroups is not known, but brute-force search for small $p$ indicates that this number is much less than the number of all subgroups of $\mathfrak{S}_p$ (see Table 1). Even though cyclic subgroups of $\mathfrak{S}_p$ are in general not 2*-closed, each cyclic group corresponds to a different coloring (see Lemma 15).
For a permutation \( \sigma \in S_p \), denote its matrix by
\[
R(\sigma) := \sum_{i=1}^{p} E_{\sigma(i)i},
\]
where \( E_{ab} \) is the \( p \times p \) matrix with 1 in the \((a, b)\)-entry and 0 in other entries. The condition \( x_{\sigma(i)\sigma(j)} = x_{ij} \) is then equivalent to \( R(\sigma) \cdot x \cdot R(\sigma)^\top = x \). Consequently,
\[
Z_\Gamma = \left\{ x \in \text{Sym}(p; \mathbb{R}) ; R(\sigma) \cdot x \cdot R(\sigma)^\top = x \text{ for all } \sigma \in \Gamma \right\}.
\]

**Definition 3.** Let \( \pi_\Gamma : \text{Sym}(p; \mathbb{R}) \to Z_\Gamma \) be the orthogonal projection on \( Z_\Gamma \), i.e. the linear map such that for any \( x \in \text{Sym}(p; \mathbb{R}) \) the element \( \pi_\Gamma(x) \in Z_\Gamma \) is uniquely determined by
\[
\text{Tr} [x \cdot y] = \text{Tr} [\pi_\Gamma(x) \cdot y] \quad (y \in Z_\Gamma).
\]

In view of (4), it is clear that
\[
\pi_\Gamma(x) = \frac{1}{|\Gamma|} \sum_{\sigma \in \Gamma} R(\sigma) \cdot x \cdot R(\sigma)^\top
\]
satisfies the above definition. Here \( |\Gamma| \) denotes the order of \( \Gamma \).

2.2. \( Z_\Gamma \) as a Jordan algebra. To derive analytic expression for Gamma-like functions on \( P_\Gamma \) it is convenient to see \( Z_\Gamma \) as a Euclidean Jordan algebra and \( P_\Gamma \) as the corresponding symmetric cone. This fact was already observed in Jensen (1988). We recall here the fundamentals of Jordan algebras, cf. Faraut and Korányi (1994). A Euclidean Jordan algebra is a Euclidean space \( A \) (endowed with the scalar product denoted by \( \langle \cdot, \cdot \rangle \)) equipped with a bilinear mapping (product)
\[
A \times A \ni (x, y) \mapsto x \bullet y \in A
\]
such that for all \( x, y, z \) in \( A \):
(i) \( x \bullet y = y \bullet x \),
(ii) \( x \bullet ((x \bullet x) \bullet y) = (x \bullet x) \bullet (x \bullet y) \),
(iii) \( \langle x, y \bullet z \rangle = \langle x \bullet y, z \rangle \).

A Euclidean Jordan algebra is said to be simple if it is not a Cartesian product of two Euclidean Jordan algebras of positive dimensions. We have the following result.

**Proposition 2.** The Euclidean space \( Z_\Gamma \) with inner product \( \langle x, y \rangle = \text{Tr} [x \cdot y] \) and the Jordan product
\[
x \bullet y = \frac{1}{2} (x \cdot y + y \cdot x),
\]
is a Euclidean Jordan algebra. This algebra is generally non-simple.

**Proof.** Since \( Z_\Gamma \) is a subset of the Euclidean Jordan algebra \( \text{Sym}(p; \mathbb{R}) \), if it is endowed with Jordan product (7), conditions (i)–(iii) are automatically satisfied. Moreover, characterization (4) of \( Z_\Gamma \) implies that the Jordan product is closed in \( Z_\Gamma \), that is, \( R(\sigma) \cdot (x \bullet y) = (x \bullet y) \cdot R(\sigma) \) for all \( x, y \in Z_\Gamma \) and \( \sigma \in \Gamma \). The result follows. \( \square \)
Up to linear isomorphism, there are only five kinds of Euclidean simple Jordan algebras. Let $\mathbb{K}$ denote the set of either the real numbers $\mathbb{R}$, the complex ones $\mathbb{C}$ or the quaternions $\mathbb{H}$. Let us write $\text{Herm}(r; \mathbb{K})$ for the space of $r \times r$ Hermitian matrices valued in $\mathbb{K}$. Then $\text{Sym}(r; \mathbb{R}), r \geq 1, \text{Herm}(r; \mathbb{C}), r \geq 2, \text{Herm}(r; \mathbb{H}), r \geq 2$ are the first three kinds of Euclidean simple Jordan algebras and they are the only ones that will concern us. The determinant and trace in Jordan algebras $\text{Herm}(r; \mathbb{K})$ will be denoted by $\det$ and $\text{tr}$ (see (Faraut and Korányi, 1994, p. 29)) respectively, so that they can be easily distinguished from the determinant and trace in $\text{Mat}(n; \mathbb{R})$ which we denote by $\det$ and $\text{Tr}$.

To each Euclidean Jordan algebra $\mathcal{A}$, one can attach the set $\Omega$ of Jordan squares, that is, $\Omega = \{ x \odot x; x \in \mathcal{A} \}$. The interior $\Omega$ of $\Omega$ is a symmetric cone, that is, it is self-dual and homogeneous. We say that $\Omega$ is irreducible if it is not the Cartesian product of two convex cones. One can prove that an open convex cone is symmetric and irreducible if and only if it is the symmetric cone $\Omega$ of some Euclidean simple Jordan algebra. Each simple Jordan algebra corresponds to a symmetric cone. The first three kinds of irreducible symmetric cones are thus, the symmetric positive definite real matrices $\text{Sym}^+(r; \mathbb{R})$ for $r \geq 1$, complex Hermitian positive definite matrices $\text{Herm}^+(r; \mathbb{C})$, and quaternionic Hermitian positive definite matrices $\text{Herm}^+(r; \mathbb{H}), r \geq 2$.

It follows from Definition 2 and Proposition 2 that $\mathcal{P}_\Gamma$ is a symmetric cone. In (Faraut and Korányi, 1994, Proposition III.4.5) it is stated that any symmetric cone is a direct sum of irreducible symmetric cones. As it will turn out, only three out of the five kinds of irreducible symmetric cones may appear in this decomposition.

Moreover, we will want to represent the elements of the symmetric cones in their real symmetric matrix representations. So, we recall that both $\text{Herm}(r; \mathbb{C})$ and $\text{Herm}(r; \mathbb{H})$ can be realized as real symmetric matrices, but of bigger dimension. For $z = a + b i \in \mathbb{C}$ define $M_C(z) = \begin{pmatrix} a - b & c - di \\ c + di & b \end{pmatrix}$. The function $M_C$ is a matrix representation of $\mathbb{C}$. Similarly, any $r \times r$ complex matrix can be realized as a $(2r) \times (2r)$ real matrix by setting the correspondence

$$\text{Mat}(r, r; \mathbb{C}) \ni (z_{i,j})_{1 \leq i,j \leq r} \simeq (M_C(z_{i,j}))_{1 \leq i,j \leq r} \in \text{Mat}(2r, 2r; \mathbb{R}),$$

that is, an $(i,j)$-entry of a complex matrix is replaced by its $2 \times 2$ real matrix representation. Note that $M_C$ maps the space $\text{Herm}(r; \mathbb{C})$ of Hermitian matrices into the space $\text{Sym}(2r; \mathbb{R})$ of symmetric matrices. For example,

$$M_C \begin{pmatrix} a & c - di \\ c + di & b \end{pmatrix} = \begin{pmatrix} a & 0 & c & d \\ 0 & a & -d & c \\ c & -d & b & 0 \\ d & c & 0 & b \end{pmatrix}. $$

Moreover, by direct calculation one sees that

$$\det \begin{pmatrix} a & 0 & c & d \\ 0 & a & -d & c \\ c & -d & b & 0 \\ d & c & 0 & b \end{pmatrix} = \det \begin{pmatrix} a & c - di \\ c + di & b \end{pmatrix}^2.$$

It can be shown that, in general,

\begin{equation}
\det (M_C(Z)) = [\det (Z)]^2 \quad \text{and} \quad \text{Tr} [M_C(Z)] = 2 \text{tr} [Z] \quad (Z \in \text{Herm}(r; \mathbb{C})).
\end{equation}

Similarly, quaternions can be realized as a $4 \times 4$ matrix:

$$a + bi + cj + dk \simeq \begin{pmatrix} a + bi & -c + di \\ c + di & a - bi \end{pmatrix} \simeq \begin{pmatrix} a - b & c - d \\ b & a & d - c \\ c - d & a & b \\ d & c & -b & a \end{pmatrix}. $$
Then, quaternionic \( r \times r \) matrices are realized as \((4r) \times (4r)\) real matrices. Thus, \( M_\mathbb{H} \) maps \( \text{Herm}(r; \mathbb{H}) \) into \( \text{Sym}(4r; \mathbb{R}) \). Moreover, it is true that

\[
\text{Det} (M_\mathbb{H}(Z)) = [\text{det} (Z)]^4 \quad \text{and} \quad \text{Tr} [M_\mathbb{H}(Z)] = 4 \text{tr}[Z] \quad (Z \in \text{Herm}(r; \mathbb{H})).
\]

2.3. Determining the structure constants and invariant measure on \( \mathcal{P}_\Gamma \). As mentioned in the introduction, in order to derive the analytic expression of the Gamma-like functions on \( \mathcal{P}_\Gamma \), we need the structure constants \((k_i, d_i, r_i)_{i=1}^L\) as well as the invariant measure \( \varphi_\Gamma \). However, due to Proposition 7 below, \( \varphi_\Gamma(X) \) is expressed in terms of the polynomials \( \text{det}(x_i) \), where \( x_i \in \text{Herm}(r_i; \mathbb{K}_i), i = 1, \ldots, L \), coming from decomposition (1). These can be derived from the decomposition of \( Z_\Gamma \). Let us note that the constants \((d_i)_{i} \) and \((k_i)_{i} \) depend only on the group \( \Gamma \), while \( r_i \) depend on a particular representation of \( \Gamma \), which is \( R \).

In view of decomposition (1), for \( X \in Z_\Gamma \), define \( \phi_i(X) = x_i \in \text{Herm}(r_i; \mathbb{K}_i) \) for \( i = 1, \ldots, L \).

**Corollary 3.** For \( X \in Z_\Gamma \), one has

\[
\text{Det}(X) = \prod_{i=1}^L \text{det}(\phi_i(X))^{k_i}.
\]

**Proof.** By (1), we have

\[
\text{Det}(X) = \prod_{i=1}^L \text{Det} (M_{\mathbb{K}_i}(x_i) \otimes I_{k_i/d_i}) = \prod_{i=1}^L \text{Det} (M_{\mathbb{K}_i}(x_i))^{k_i/d_i} = \prod_{i=1}^L \left[ \text{det}(x_i)^{d_i} \right]^{k_i/d_i} = \prod_{i=1}^L \text{det}(x_i)^{k_i},
\]

whence follows the formula. We have used (8) and (9) for the third equality above. \( \square \)

**Lemma 4.** Assume that \( \Gamma \subset \mathfrak{S}_p \) and that \((k_i, d_i, r_i)_{i=1}^L\) are the structure constants corresponding to \( Z_\Gamma \). Assume that we have an irreducible factorization

\[
\text{Det}(X) = \prod_{j=1}^{L'} f_j(X)^{a_j} \quad (X \in Z_\Gamma),
\]

where each \( a_j \) is a positive integer, each \( f_j(X) \) is an irreducible polynomial of \( X \in Z_\Gamma \), and \( f_i \neq f_j \) if \( i \neq j \).

Then, \( L = L' \), for each \( j \) there exists unique \( i \) such that \( f_j(X)^{a_j} = \text{det}(\phi_i(X))^{k_i} \), and

a) \( k_i = a_j \),

b) \( r_i \) is the degree of \( f_j(X) = \text{det}(\phi_i(X)) \),

c) if \( r_i > 1 \), then \( d_i \) can be calculated from \( r_i + d_i(r_i - 1)/2 = \text{rank}(P_j) \), where \( P_j \) is the linear operator defined by \( Z_\Gamma \ni x \mapsto P_j(x) = E_j \cdot x \in Z_\Gamma \) and \( E_j \in Z_\Gamma \) is the gradient of \( f_j(X) \) at \( X = I_p \).

**Remark 4.** If \( r_i = 1 \), the determination of \( d_i \) is not needed for writing the block decomposition of \( Z_\Gamma \), since in this case \( \mathbb{R} = \text{Herm}(1; \mathbb{R}) = \text{Herm}(1; \mathbb{C}) = \text{Herm}(1; \mathbb{H}) \) and, if \( k_i \) is divisible by 2 or by 4, we have \( M_{\mathbb{K}_i}(x_i) \otimes I_{k_i/d_i} = x_i I_{k_i} \).
Proof of Lemma 4. Since the determinant polynomial of a simple Jordan algebra is always irreducible (Upmeier, 1986, Lemma 2.3 (1)), comparing (10) and (11), we obtain \( L = L' \), and that, for each \( j \), there exists \( i \) such that \( f_j(X)^{d_i} = \det(\phi_i(X))^{k_i} \). From this follows also \( a \) and \( b \).

Observe that \( r_i + d_i r_i (r_i - 1)/2 = \dim_{\mathbb{R}} \text{Herm}(r_i; \mathbb{K}_i) \). Point \( c \) follows from the fact that \( P_j \) coincides with the projection \( \bigoplus_{i=1}^L \text{Herm}(r_i; \mathbb{K}_i) \rightarrow \text{Herm}(r_i; \mathbb{K}_i) \).

The practical significance of the method proposed in this lemma is that neither representation theory nor group theory is used. It is a strong advantage when we consider colorings corresponding to a large number of different groups, for which finding structure constants is very complicated.

Remark 5. The factorization of multivariate polynomials over an algebraic number field can be done for example in PYTHON (see sympy.polys.polytools.factor) or in MATHEMATICA (see Factor). However, in order to make use of Lemma 4, one has to perform a factorization over the real number field. It turns out that the previously listed tools can be used for this purpose by selecting an appropriate Extension parameter. Indeed, in our setting, the irreducible factorization over the real number field coincides with the one over the real cyclotomic field

\[
\mathbb{Q} \left[ \zeta + \frac{1}{\zeta} \right] = \left\{ \sum_{k=0}^{\varphi_E(M)/2-1} q_k \left( \zeta + \frac{1}{\zeta} \right)^k ; q_k \in \mathbb{Q}, k = 0,1,\ldots,\varphi_E(M)/2-1 \right\},
\]

where \( \zeta \) is the primitive \( M \)-th root \( e^{2\pi i/M} \) of unity with \( M \) being the least common multiple of the orders of elements \( \sigma \in \Gamma \), and \( \varphi_E(M) \) is the number of positive integers up to \( M \) that are relatively prime to \( M \) (Serre, 1977, Section 12.3).

An example showing the utility of Lemma 4 can be found in the Supplementary material (Graczyk et al., Section 3.1).

2.4. Finding structure constants and construction of the orthogonal matrix \( U_\Gamma \) when \( \Gamma \) is cyclic. We now show that, when the group \( \Gamma \) is generated by one permutation \( \sigma \in \mathfrak{S}_p \), the orthogonal matrix \( U_\Gamma \) can be constructed explicitly, and we obtain the structure constants \( r_i, k_i \) and \( d_i \) easily.

Let us consider the \( \Gamma \)-orbits in \( \{1,2,\ldots,p\} \). Let \( \{i_1,\ldots,i_C\} \) be a complete system of representatives of the \( \Gamma \)-orbits, and for each \( c = 1,\ldots,C \), let \( p_c \) be the cardinality of the \( \Gamma \)-orbit through \( i_c \). The order \( N \) of \( \Gamma \) equals the least common multiple of \( p_1, p_2,\ldots,p_C \) and one has \( \Gamma = \{\text{id}, \sigma, \sigma^2,\ldots,\sigma^{N-1}\} \). In what follows, we treat 0 as a multiple of \( N \).

Theorem 5. Let \( \Gamma = \langle \sigma \rangle \) be a cyclic group of order \( N \). For \( \alpha = 0,1,\ldots,\left\lfloor \frac{N}{2} \right\rfloor \) set

\[
r^*_\alpha = \# \{ c \in \{1,\ldots,C\}; \alpha p_c \text{ is a multiple of } N \}
\]

\[
d^*_\alpha = \begin{cases} 1 & (\alpha = 0 \text{ or } N/2) \\ 2 & \text{(otherwise).} \end{cases}
\]

Then we have \( L = \# \{ \alpha ; r^*_\alpha > 0 \} \), \( r = (r^*_\alpha ; r^*_\alpha > 0) \) and \( k = d = (d^*_\alpha ; r^*_\alpha > 0) \).

Note that, \( r^*_0 \) equals the number \( C \) of cycles in a decomposition of a permutation.
EXAMPLE 6. Let us consider $\sigma = (1, 2, 3)(4, 5)(6) \in S_6$. The three $\Gamma$-orbits are \{1, 2, 3\}, \{4, 5\} and \{6\}. Set $i_1 = 1$, $i_2 = 4$, $i_3 = 6$. Then $p_1 = 3$, $p_2 = 2$, $p_3 = 1$. We have $N = 6$. We count $r_0^* = 3$, $r_1^* = 0$, $r_2^* = 1$, $r_3^* = 1$, so that $r = (3, 1, 1)$. Since $d = (1, 2, 1)$, we have $Z_\Gamma \simeq \text{Sym}(3; \mathbb{R}) \oplus \text{Herm}(1; \mathbb{C}) \oplus \text{Sym}(1; \mathbb{R})$.

For $c = 1, \ldots, C$ define $v_1^{(c)}(\cdot), \ldots, v_{p_c}^{(c)} \in \mathbb{R}^p$ by

$$v_1^{(c)} := \sqrt{\frac{1}{p_c}} \sum_{k=0}^{p_c-1} e_{\sigma^k(i_c)},$$

$$v_{2\beta}^{(c)} := \sqrt{\frac{2}{p_c}} \sum_{k=0}^{p_c-1} \cos\left(\frac{2\pi \beta k}{p_c}\right) e_{\sigma^k(i_c)} \quad (1 \leq \beta < p_c/2),$$

$$v_{2\beta+1}^{(c)} := \sqrt{\frac{2}{p_c}} \sum_{k=0}^{p_c-1} \sin\left(\frac{2\pi \beta k}{p_c}\right) e_{\sigma^k(i_c)} \quad (1 \leq \beta < p_c/2),$$

$$v_{p_c}^{(c)} := \sqrt{\frac{1}{p_c}} \sum_{k=0}^{p_c-1} \cos(\pi k) e_{\sigma^k(i_c)} \quad (\text{if } p_c \text{ is even}).$$

THEOREM 6. The orthogonal matrix $U_\Gamma$ from Theorem 1 can be obtained by arranging column vectors $\{v_k^{(c)}\}_{1 \leq c \leq C, 1 \leq k \leq p_c}$ in the following way: we put $v_k^{(c)}$ earlier than $v_k^{(c')}$ if

(i) $\frac{k}{[k/2]} < \frac{k'}{[k'/2]}$, or

(ii) $\frac{p_c}{[k/2]} = \frac{p_c'}{[k'/2]}$ and $c < c'$, or

(iii) $\frac{p_c}{[k/2]} = \frac{p_c'}{[k'/2]}$ and $c = c'$ and $k$ is even and $k'$ is odd.

Proofs of the above results are presented in the Supplementary material. We shall see there that $R(\sigma)$ acts on the 2-dimensional space spanned by $v_{2\beta}^{(c)}$ and $v_{2\beta+1}^{(c)}$ as a rotation with the angle $2\pi \beta / p_c, 1 \leq \beta < p_c/2$. The condition (i) means that the angle for $v_k^{(c)}$ is smaller than the one for $v_k^{(c')}$.  

EXAMPLE 7. We continue Example 6. According to Theorem 6,

$$U_\Gamma = \begin{pmatrix} v_1^{(1)}, v_1^{(2)}, v_1^{(3)}, v_2^{(1)}, v_2^{(2)} \end{pmatrix} = \begin{pmatrix} 1/\sqrt{3} & 0 & 0 & \sqrt{2/3} & 0 & 0 \\ 1/\sqrt{3} & 0 & 0 & -\sqrt{1/6} & 1/\sqrt{2} & 0 \\ 1/\sqrt{3} & 0 & 0 & -\sqrt{1/6} & -1/\sqrt{2} & 0 \\ 0 & 1/\sqrt{2} & 0 & 0 & 0 & 1/\sqrt{2} \\ 0 & 1/\sqrt{2} & 0 & 0 & -1/\sqrt{2} & 0 \\ 0 & 0 & 1 & 0 & 0 & 0 \end{pmatrix}.$$  

Then we have

$$U_\Gamma^T \cdot R(\sigma^k) \cdot U_\Gamma = \begin{pmatrix} I_3 \otimes B_0(\sigma^k) \\ B_2(\sigma^k) \\ B_3(\sigma^k) \end{pmatrix},$$

where $B_0(\sigma^k) = 1$, $B_2(\sigma^k) = \text{Rot}(\frac{2\pi k}{3}) \in \text{GL}(2; \mathbb{R})$ and $B_3(\sigma^k) = (-1)^k$. Here $\text{Rot}(\theta)$ denotes the rotation matrix $\begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}$ for $\theta \in \mathbb{R}$.  

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The block diagonal decomposition of $Z_{\Gamma}$ is

$$U_{\Gamma}^{\top} \cdot Z_{\Gamma} \cdot U_{\Gamma} = \left\{ \begin{pmatrix} x_1 \\ x_2 I_2 \\ x_3 \end{pmatrix} ; x_1 \in \text{Sym}(3; \mathbb{R}), \ x_2, x_3 \in \mathbb{R} \right\}.$$ 

**Remark 8.** In the cyclic case we have $k = d$ and so the formula (1) holds without the Kronecker product terms. Since $d_i \in \{1, 2\}$, the quaternionic case never occurs.

3. **Gamma integrals and normalizing constants.**

3.1. **Gamma integrals on irreducible symmetric cones.** Let $\Omega$ be one of the first three kinds of irreducible symmetric cones, that is, $\Omega = \text{Herm}^+(r; \mathbb{K})$, where $\mathbb{K} \in \{\mathbb{R}, \mathbb{C}, \mathbb{H}\}$. As before, determinant and trace on corresponding Euclidean Jordan algebras are denoted by $\det$ and $\tr$. Then, we have the relation

$$\dim \Omega = r + \frac{r(r-1)}{2}d,$$

where $d = 1$ if $\mathbb{K} = \mathbb{R}$, $d = 2$ if $\mathbb{K} = \mathbb{C}$ and $d = 4$ if $\mathbb{K} = \mathbb{H}$.

Recall that Euclidean measure is the volume measure induced by the Euclidean metric. Let $m(dx)$ denote the Euclidean measure associated with the Euclidean structure defined on $\mathcal{A} = \text{Herm}(r; \mathbb{K})$ by $\langle x, y \rangle = \tr[x \cdot y] = \tr[x \cdot y]$. The Gamma integral

$$\Gamma_{\Omega}(\lambda) := \int_{\Omega} \det(x)^{\lambda} e^{-\tr[x]} \det(x)^{-\dim \Omega/r} m(dx)$$

is finite if and only if $\lambda > \frac{1}{2}(r-1)d = \dim \Omega/r - 1$ and in such case

$$\Gamma_{\Omega}(\lambda) = (2\pi)^{(\dim \Omega - r)/2} \Gamma(\lambda) \Gamma(\lambda - d/2) \ldots \Gamma(\lambda - (r-1)d/2).$$

Moreover, one has

$$\int_{\Omega} \det(x)^{\lambda} e^{-\tr[x-y]} \det(x)^{-\dim \Omega/r} m(dx) = \Gamma_{\Omega}(\lambda) \det(y)^{-\lambda}$$

for any $y \in \Omega$.

The measure $\mu_{\Omega}(dx) = \det(x)^{-\dim \Omega/r} m(dx)$ is invariant in the following sense. Let $G(\Omega)$ be the linear automorphism group of $\Omega$, that is, the set $\{ g \in \text{GL}(\mathcal{A}) ; g \Omega = \Omega \}$, where $\mathcal{A}$ is the associated Euclidean Jordan algebra. Then, the measure $\mu_{\Omega}$ is a $G(\Omega)$-invariant measure in the sense that for any Borel measurable set $B$ one has

$$\mu_{\Omega}(g^{-1}B) = \mu_{\Omega}(B) \quad (g \in G(\Omega)).$$

3.2. **Gamma integrals on the cone $P_{\Gamma}$.** We endow the space $Z_{\Gamma}$ with the scalar product

$$\langle x, y \rangle = \text{Tr}[x \cdot y] \quad (x, y \in Z_{\Gamma}).$$

Let $dX$ denote the Euclidean measure on the Euclidean space $(Z_{\Gamma}, \langle \cdot, \cdot \rangle)$. Let us note that this normalization is not important in the Bayesian model selection procedure as there we always consider quotients of integrals.

**Example 9.** Consider $p = 3$ and $\Gamma = S_3$. The space $Z_{\Gamma}$ is 2-dimensional and it consists of matrices of the form (see Example 1)

$$X = \begin{pmatrix} a & b & b \\ b & a & b \\ b & b & a \end{pmatrix}$$
for \( a, b \in \mathbb{R} \). Since \( \|X\|^2 = \text{Tr} [X^2] = 3a^2 + 6b^2 = v^T v \) with \( v^T = (\sqrt{3}a, \sqrt{6}b) \), we have \( dX = \sqrt{3}\sqrt{6} da \, db = 3\sqrt{2} da \, db \).

Generally, if \( m_i \) denotes the Euclidean measure on \( \mathcal{A}_i := \text{Herm}(r_i; \mathbb{K}_i) \) with the inner product defined from the Jordan algebra trace (recall (8) and (9)), then (1) implies that for \( X \in \mathcal{Z}_\Gamma \) we have

\[
\|X\|^2 = \langle X, X \rangle = \sum_{i=1}^{L} \frac{k_i}{d_i} \text{Tr} \left[ M_{\mathcal{K}_i}(x_i)^2 \right] = \sum_{i=1}^{L} k_i \text{tr}[(x_i)^2],
\]

which implies that

\[
dX = \prod_{i=1}^{L} \left( \sqrt{k_i} \right)^{\dim \Omega_i} \, m_i(dx_i) = e^{B_{\Gamma}} \prod_{i=1}^{L} m_i(dx_i),
\]

where

\[
B_{\Gamma} := \frac{1}{2} \sum_{i=1}^{L} (\dim \Omega_i)(\log k_i).
\]

**Definition 10.** Let \( G(\mathcal{P}_\Gamma) = \{ g \in \text{GL}(p; \mathbb{R}) \mid g\mathcal{P}_\Gamma = \mathcal{P}_\Gamma \} \) be the linear automorphism group of \( \mathcal{P}_\Gamma \). We define the \( G(\mathcal{P}_\Gamma) \)-invariant measure \( \varphi_{\Gamma}(X) \, dX \) by

\[
\varphi_{\Gamma}(X) = e^{B_{\Gamma}} \left( \prod_{i=1}^{L} \frac{1}{\Gamma_{\Omega_i}(\dim \Omega_i/r_i)} \right) \int_{\mathcal{P}_\Gamma^*} e^{-\text{Tr}[X \cdot Z]} \, dZ,
\]

where \( \mathcal{P}_\Gamma^* = \{ y \in \mathcal{Z}_\Gamma \mid \text{Tr} [y \cdot x] > 0, \forall x \in \mathcal{P}_\Gamma \setminus \{0\} \} \) is the dual cone of \( \mathcal{P}_\Gamma \).

**Proposition 7.** We have

\[
\varphi_{\Gamma}(X) = \prod_{i=1}^{L} \left( \det(\phi_i(X)) \right)^{-\dim \Omega_i/r_i}.
\]

The proofs of the following results of this section can be found in the supplementary material Graczyk et al.

**Definition 11.** The Gamma function of \( \mathcal{P}_\Gamma \) is defined by the following integral

\[
\Gamma_{\mathcal{P}_\Gamma}(\lambda) := \int_{\mathcal{P}_\Gamma} \text{Det} (X)^\lambda e^{-\text{Tr}[X]} \varphi_{\Gamma}(X) \, dX,
\]

whenever it converges.

**Theorem 8.** The integral (17) converges if and only if

\[
\lambda > \max_{i=1,\ldots,L} \left\{ \frac{(r_i - 1)d_i}{2k_i} \right\}
\]

and, for these values of \( \lambda \), we have

\[
\Gamma_{\mathcal{P}_\Gamma}(\lambda) = e^{-A_{\Gamma} \lambda + B_{\Gamma}} \prod_{i=1}^{L} \Gamma_{\Omega_i}(k_i \lambda)
\]
where $\Gamma_\Omega_i$ is given in (12), $B_\Gamma$ in (15) and

\begin{equation}
A_\Gamma := \sum_{i=1}^L r_i k_i \log k_i.
\end{equation}

Moreover, if $Y \in \mathcal{P}_\Gamma$ and (18) holds true, then

\begin{equation}
\int_{\mathcal{P}_\Gamma} \det(X)^\lambda e^{-\text{Tr}[Y \cdot X]} \varphi_\Gamma(X) \, dX = \Gamma_{\mathcal{P}_\Gamma}(\lambda) \det(Y)^{-\lambda}.
\end{equation}

We also have the following result

**Theorem 9.** If $Y \in \mathcal{P}_\Gamma$ and

\[ \lambda > \max_{i=1, \ldots, L} \left\{ -\frac{1}{k_i} \right\}, \]

then

\begin{equation}
\int_{\mathcal{P}_\Gamma} \det(X)^\lambda e^{-\text{Tr}[Y \cdot X]} \, dX = e^{-A_\Gamma \lambda - B_\Gamma} \prod_{i=1}^L \Gamma_{\Omega_i} \left( k_i \lambda + \frac{\dim \Omega_i}{r_i} \right) \frac{\varphi_\Gamma(Y)}{\det(Y)^\lambda}.
\end{equation}

3.3. **RCOP-Wishart laws on $\mathcal{P}_\Gamma$.** Let $\Sigma \in \mathcal{P}_\Gamma \subset \text{Sym}^+(p; \mathbb{R})$ and consider i.i.d. random vectors $Z^{(1)}, \ldots, Z^{(n)}$ following the $\mathcal{N}(0, \Sigma)$ distribution. Define $U_i = Z^{(i)} \cdot Z^{(i)\top}$, $i = 1, \ldots, n$, and $U = \sum_{i=1}^n U_i$. We note that such model is clearly not invariant under changing the scale of variables: random vector $\text{diag}(\alpha) \cdot Z^{(1)}$ for $\alpha \in \mathbb{R}^p$ is in general not invariant under any permutation subgroup. Such issue is an immanent property of RCON models (a generalization of RCOP models) and was noticed already in Højsgaard and Lauritzen (2008). The authors recommend to keep all variables in the same units.

Our aim is to analyze the probability distribution of the random matrix

\[ W_n = \pi_\Gamma(U) = \pi_\Gamma(U_1 + \cdots + U_n) = \pi_\Gamma(U_1) + \ldots + \pi_\Gamma(U_n). \]

In the rest of this section, we find $n_0$ such that for $n \geq n_0$ the random matrix $W_n$ follows an absolutely continuous law, and we compute its density. Further, we extend the shape parameter to a continuous range and define the RCOP-Wishart law on $\mathcal{P}_\Gamma$.

We start with the following easy result.

**Lemma 10.** For any $\theta \in \text{Sym}^+(p; \mathbb{R})$ we have

\[ \mathbb{E}e^{-\text{Tr}[\theta \cdot \pi_\Gamma(U_1)]} = \det(I_p + 2 \Sigma \cdot \pi_\Gamma(\theta))^{-1/2}. \]

**Proof.** Using (5) repeatedly we have

\[ \text{Tr} [\theta \cdot \pi_\Gamma(U_1)] = \text{Tr}[\pi_\Gamma(\theta) \cdot \pi_\Gamma(U_1)] = \text{Tr}[\pi_\Gamma(\theta) \cdot U_1]. \]

The assertion follows from the usual multivariate Gauss integral. \qed

**Proposition 11.** The law of $W_n$ is absolutely continuous on $\mathcal{P}_\Gamma$ if and only if

\[ n \geq n_0 := \max_{i=1, \ldots, L} \left\{ \frac{r_i d_i}{k_i} \right\}. \]

If $n \geq n_0$, then its density function with respect to $dX$ is given by

\begin{equation}
\frac{\det(X)^{n/2} e^{-\frac{1}{2} \text{Tr}[X \cdot \Sigma^{-1}]} \varphi_\Gamma(X)}{\det(2\Sigma)^{n/2} \Gamma_{\mathcal{P}_\Gamma}(\frac{n}{2})} \varphi_\Gamma(X) \mathbb{1}_{\mathcal{P}_\Gamma}(X).
\end{equation}
PROOF. With $\lambda = n/2$, condition (18) becomes
\[ n > \max_{i=1, \ldots, L} \left\{ \frac{(r_i - 1)d_i}{k_i} \right\}. \]
Since the quotient $k_i/d_i$ is an integer, the last condition is equivalent to (23).

In view of Lemma 10, it is enough to show that $W_n$ has density (24) if and only if for any $\theta \in \mathcal{P}_\Gamma$,
\[ \mathbb{E} e^{-\text{Tr}[\theta \cdot W_n]} = \text{Det} (I_p + 2 \Sigma \cdot \theta)^{-n/2}. \] (25)
This follows directly from (21).

It is known that the MLE exists and is unique if and only if the sufficient statistic lies in the interior of its convex support, see Barndorff-Nielsen (2014). It is clear that if (23) is not satisfied, then the support of $W_n$ is contained in the boundary of $\mathcal{P}_\Gamma$. Recall that the orthogonal projection $\pi_\Gamma$ is given by (6).

COROLLARY 12. The MLE of $\Sigma$ exists if and only if the number of samples $n$ satisfies (23). If it exists, it is given by
\[ \hat{\Sigma} = \frac{1}{n} \pi_\Gamma (U_1 + \cdots + U_n). \]

The above result has been already proven in (Andersson, 1975, Theorem 5.9) (see also (Andersson and Madsen, 1998, Sec. A.3, A.4)).

REMARK 12. If $U = (U_1 + \cdots + U_n)/n$ is positive definite, then $U$ can be regarded as an empirical covariance matrix. The Kullback-Leibler divergence between $N_p(0, U)$ and $N_p(0, \Sigma)$ is equal to \( \frac{1}{2} \left\{ \log \text{det} \Sigma + \text{tr} U \Sigma^{-1} - \log \text{det}(U) - p \right\} \), which is obviously minimized by the MLE $\hat{\Sigma}$. Therefore, Corollary 12 implies that $\pi_\Gamma(U)$ is the Kullback-Leibler projection of $U$ onto $\mathcal{P}_\Gamma$. Goutis and Robert (1998). We note that the KL projection in general is not linear, whereas $\pi_\Gamma$ clearly is.

Let us recall that the MLE of $\Sigma$ in the standard normal model exists if and only if $n \geq p$. We recover this case for $\Gamma = \{ \text{id} \}$, since then we have $L = 1$, $r_1 = p$ and $k_1 = d_1 = 1$.

When $n < n_0$, the law of $W_n$ is singular, and it can be described as a direct product of the singular Wishart laws on the irreducible symmetric cones $\Omega_i$, see e.g. Hassairi and Lajmi (2001).

DEFINITION 13. Let $\eta > \max \left\{ (r_i - 1)d_i/k_i ; i = 1, \ldots, L \right\}$ and $\Sigma \in \mathcal{P}_\Gamma$. The RCOP–Wishart law $W_{\Gamma, \Sigma}^\eta$ is defined by its density
\[ W_{\eta, \Sigma}^\Gamma(dX) = \frac{\text{Det}(X)^{\eta/2} e^{-\frac{1}{2} \text{Tr}[X^{-1} \Sigma^{-1}]} }{\text{Det}(2\Sigma)^{\eta/2} \Gamma_{\text{Pr}}(\frac{\eta}{2})} \varphi_\Gamma(X) 1_{\mathcal{P}_\Gamma}(X) dX. \] (26)
With this new notation, we see that if (23) is satisfied, then $W_n \sim W_{\Gamma, \Sigma}^\eta$.

LEMMA 13. The Jacobian of the transformation $\mathcal{P}_\Gamma \ni X \mapsto X^{-1} \in \mathcal{P}_\Gamma$ equals $\varphi_\Gamma(X^{-1})^2$. 

Proof of the lemma can be found in the Supplementary material. By this lemma we obtain another useful formula for the invariant measure, namely

$$\varphi_{\Gamma}(X) = \text{Det}_{\text{End}}(\mathbb{P}_X)^{-1/2}~~(X \in \mathbb{Z}_\Gamma),$$

where \(\text{Det}_{\text{End}}\) is the determinant in the space of endomorphisms of \(\mathbb{Z}_\Gamma\) and for any \(X \in \mathbb{Z}_\Gamma\) by \(\mathbb{P}_X\) we denote the linear map on \(\mathbb{Z}_\Gamma\) to itself defined by \(\mathbb{P}_X Y = X \cdot Y \cdot X\). Lemma 13 gives also the following result.

**Proposition 14.** Let \(W \sim W_{\eta,\Sigma}^\Gamma\) with \(\eta > \max \{ (r_i - 1)d_i/k_i; i = 1, \ldots, L \}\) and \(\Sigma \in \mathcal{P}_\Gamma\). Then its inverse \(Y = W^{-1}\) has density

$$\frac{\text{Det}(Y)^{-\eta/2}e^{-\frac{1}{2} \text{Tr}[Y^{-1}\Sigma^{-1}]}}{\text{Det}(2\Sigma)^{\eta/2} \Gamma_{\mathcal{P}_\Gamma}(\frac{\eta}{2})} \varphi_{\Gamma}(Y)1_{\mathcal{P}_\Gamma}(Y).$$

**3.4. The Diaconis-Ylvisaker conjugate prior for \(K\).** The Diaconis-Ylvisaker conjugate prior (Diaconis and Ylvisaker (1979)) for the canonical parameter \(K = \Sigma^{-1}\) is given by

$$f(K; \delta, D) = \frac{1}{I_{\Gamma}(\delta, D)} \text{Det}(K)^{(\delta-2)/2}e^{-\frac{1}{2} \text{Tr}[K^{-1}D]}1_{\mathcal{P}_\Gamma}(K),$$

for hyper-parameters \(\delta > 2 \max \{ 1 - 1/k_i; i = 1, \ldots, L \}\) and \(D \in \mathcal{P}_\Gamma\). By (22), the normalizing constant is equal to

$$I_{\Gamma}(\delta, D) = e^{-A_{\Gamma}(\delta-2)/2-B_{\Gamma}} \prod_{i=1}^{L} \Gamma_{\Omega_i}\left(\frac{\delta - 2}{2} + \frac{\dim \Omega_i}{r_i}\right) \frac{\varphi_{\Gamma}(\frac{1}{2}D)}{\text{Det}(\frac{1}{2}D)^{(\delta-2)/2}},$$

where \(A_{\Gamma}, B_{\Gamma}\) and \(\varphi_{\Gamma}\) are given in (20), (15) and (16).

We note that despite the fact that the choice of hyper-parameters is not scale invariant, statisticians usually take \(\delta = 3\) and \(D = I_p\), see e.g. Massam, Li and Gao (2018).

**4. Model selection.** Bayesian model selection on all colored spaces seems at the moment intractable. This is due in great part to a poor combinatorial description of the colored spaces \(\mathbb{Z}_\Gamma\). In particular, the number of such spaces, that is, \(# \{ \mathbb{Z}_\Gamma; \Gamma \in \mathfrak{S}_p \}\) is generally unknown for large \(p\). It was shown in Gehrmann (2011) that these colorings constitute a lattice with respect to the usual inclusion of subspaces. However the structure of this lattice is rather complicated and is unobtainable for big \(p\). This, in turn, does not allow to define a Markov chain with known transition probabilities on such colorings. Finally, the fundamental problem which prevents us from doing Bayesian model selection on all colored spaces for arbitrary \(p\) is the following. In order to compute Bayes factors, one has to be able to find the structure constants \((k_i, d_i, r_i)_{i=1}^{L}\) for arbitrary subgroups of \(\mathfrak{S}_p\). This is equivalent to finding irreducible representations over reals for an arbitrary finite group, which is very hard in general, although general algorithms have been developed for this issue (see Plesken and Souvignier (1996)).

In this section, we are making a step forward in the problem of model selection for colored models in two ways. In Section 4.1, we use the results of Section 2.4, to obtain the structure constants when we restrict our search to the space of colored models generated by a cyclic group, that is, when \(\Gamma = \langle \sigma \rangle\) for \(\sigma \in \mathfrak{S}_p\) and we propose a model selection procedure restricted to the cyclic colorings. In Section 4.2, we use Lemma 4 and Remark 5 to obtain the irreducible representations of \(\mathbb{Z}_\Gamma\) and the structure constants by factorization of the determinant. We apply this technique to do model selection for the four-dimensional example given by Frets’ data since, in that case, there are only 22 models and we can compute all the Bayes factors.
4.1. Model selection within cyclic groups. The smaller space of cyclic colorings has a much better combinatorial description. In particular, the following result can be proved.

**Lemma 15.** If $\mathcal{Z}(\sigma) = \mathcal{Z}(\sigma')$ for some $\sigma, \sigma' \in S_p$, then $\langle \sigma \rangle = \langle \sigma' \rangle$.

This result allows us to calculate the number of different colorings corresponding to cyclic groups, that is, the number of labeled cyclic subgroups of the symmetric group $S_p$, which can be found in OEIS, sequence A051625 (see the last column of Table 1).

We will present two applications of the Metropolis-Hastings algorithm. In the first one, the Markov chain will move on the space of cyclic groups. The drawback of this first approach is that we need to compute the proposal distribution $g$, whose computational complexity grows faster than quadratically as $p$ increases (see (29)). In the second algorithm, we consider a larger state space $\mathcal{S}_p$, which allows us to consider an easy proposal distribution. However, this comes at the cost of slower convergence of the posterior probabilities (see Theorem 16).

4.1.1. First approach. Each cyclic subgroup $\Gamma$ can be uniquely represented by a permutation, which is minimal in the lexicographic order within permutations generating $\Gamma$. Let $\nu(\Gamma) \in \mathcal{S}_p$ be such a permutation, that is,

$$\nu(\Gamma) = \min \{ \sigma \in \mathcal{S}_p : \langle \sigma \rangle = \Gamma \}.$$  

Define

$$c_t := \langle \nu(c_{t-1}) \circ x_t \rangle,$$

where $c_0$ is a fixed cyclic subgroup and $(x_t)_{t \in \mathbb{N}}$ is a sequence of i.i.d. random transpositions distributed uniformly, that is, $\mathbb{P}(x_t = \alpha) = 1/\binom{p}{2}$ for any $\alpha \in T := \{(i, j) \in \mathcal{S}_p\}$. Clearly, the sequence $(c_t)_t$ is a Markov chain. Its state space is the set of all cyclic subgroups of $\mathcal{S}_p$. Moreover, the trivial subgroup $\{id\}$ can be reached from any subgroup $c_t$ (and vice versa) in a finite number of steps with positive probability. Thus the chain $(c_t)_t$ is irreducible. The proposal distribution in the Metropolis-Hastings algorithm is the conditional distribution of $c_t | c_{t-1}$. It is proportional to the number of possible transitions from $c$ to $c'$, that is,

$$g(c' | c) := \frac{\# \{ (i, j) \in \mathcal{S}_p : c' = \langle \nu(c) \circ (i, j) \rangle \}}{\binom{p}{2}},$$

where $c$ and $c'$ are cyclic subgroups.

We follow the principles of Bayesian model selection for graphical models, presented, for example, in (Maathuis et al., 2018, Chapter 10, p.247). Let $\Gamma$ be uniformly distributed on...
the set $\mathcal{C} := \{ \langle \sigma \rangle : \sigma \in \mathfrak{S}_p \}$ of cyclic subgroups of $\mathfrak{S}_p$. We assume that $K|\{ \Gamma = c \}$, $c \in \mathcal{C}$, follows the Diaconis-Ylvisaker conjugate prior distribution on $P_c$ with hyper-parameters $\delta$ and $D$, that is,

$$I_{K|\Gamma = c}(k) = \frac{1}{I_c(\delta, D)} \text{Det} \left( k \right)^{(\delta-2)/2} e^{-\frac{1}{2} \text{Tr}[D-k]} 1_{P_c}(k),$$

where the normalizing constant is given in (27). Suppose that $Z_1, \ldots, Z_n$ given $\{ K = k, \Gamma = c \}$ are i.i.d. $N_p(0, k^{-1})$ random vectors with $k \in P_c$. Then, it is easily seen that we have

$$P(\Gamma = c | Z_1, \ldots, Z_n) \propto \frac{I_c(\delta + n, D + U)}{I_c(\delta, D)} \quad (c \in \mathcal{C})$$

with $U = \sum_{i=1}^n Z_i \cdot Z_i^\top$. These derivations allow us to run the Metropolis-Hastings algorithm restricted to cyclic groups, as follows.

Algorithm 14. Starting from a cyclic group $C_0 \in \mathcal{C}$, repeat the following two steps for $t = 1, 2, \ldots$:

1. Sample $x_t$ uniformly from the set $T$ of all transpositions and set $c' = \langle \nu(C_{t-1}) \circ x_t \rangle$;
2. Accept the move $C_t = c'$ with probability

$$\min \left\{ 1, \frac{I_c(\delta + n, D + U)}{I_c(\delta, D)} \frac{g(C_{t-1}|c')}{g(c'|C_{t-1})} \right\}$$

If the move is rejected, set $C_t = C_{t-1}$.

4.1.2. Second approach. It is known that $\langle \sigma \rangle = \langle \sigma' \rangle$ if and only if $\sigma' = \sigma^k$ for some $k \in \beta(|\sigma|)$, where

$$\beta(n) = \{ k \in \{1, \ldots, n \} : k \text{ and } n \text{ are relatively prime} \}$$

and $|\sigma|$ denotes the order of $\sigma$. Let $\mathcal{C} = \{ \langle \sigma \rangle : \sigma \in \mathfrak{S}_p \}$ denote the set of cyclic subgroups of $\mathfrak{S}_p$. For $c \in \mathcal{C}$ we define $\Phi(c) := \#\beta(|c|)$ and $\mathcal{C}_c := \{ \sigma \in \mathfrak{S}_p : \langle \sigma \rangle = c \}$, the set of permutations, which generate the cyclic subgroup $c$. We have

$$\Phi(c) = \#\mathcal{C}_c \quad (c \in \mathcal{C}).$$

For $c \in \mathcal{C}$, we denote

$$\pi_c = P(\Gamma = c | Z_1, \ldots, Z_n),$$

which we want to approximate. In our model we have (see (30))

$$\pi_c \propto \frac{I_c(\delta + n, D + U)}{I_c(\delta, D)} \quad (c \in \mathcal{C}).$$

In order to find $\pi = (\pi_c; c \in \mathcal{C})$ let us consider $\tilde{\pi} = (\tilde{\pi}_\sigma; \sigma \in \mathfrak{S}_p)$, a probability distribution on $\mathfrak{S}_p$ such that

$$\tilde{\pi}_\sigma \propto \frac{I_{\langle \sigma \rangle}(\delta + n, D + U)}{I_{\langle \sigma \rangle}(\delta, D)} \quad (\sigma \in \mathfrak{S}_p).$$

Since (32) and (33) imply that $\tilde{\pi}_\sigma \propto \pi_{\langle \sigma \rangle}$, we have

$$\tilde{\pi}_\sigma = \frac{\pi_{\langle \sigma \rangle}}{\sum_{c \in \mathcal{C}} \Phi(c) \pi_c} \quad (\sigma \in \mathfrak{S}).$$
As before, let \((x_t)_{t \in \mathbb{N}}\) be a sequence of i.i.d random transpositions distributed uniformly on \(T = \{(i, j) \in \mathcal{G}_p\}\). We define a random walk on \(\mathcal{G}_p\) by
\[
s_{t+1} = s_t \circ x_{t+1}, \quad (t = 0, 1, \ldots).
\]
Then, \((s_t)_t\) is an irreducible Markov chain with symmetric transition probability
\[
g(\sigma'|\sigma) = \begin{cases} \frac{1}{t}, & \text{if } \sigma^{-1} \circ \sigma' \in T, \\ 0, & \text{if } \sigma^{-1} \circ \sigma' \notin T. \end{cases}
\]
We note that \((\langle s_t \rangle)_t\) is not a Markov chain on the space of cyclic subgroups. Indeed, it can be shown that the necessary conditions for \((f(s_t))_t\) to be a Markov chain (see (Burke and Rosenblatt, 1958, Eq. (3))) are not satisfied for \(f(\sigma) := \langle \sigma \rangle\) if \(p > 4\). A remedy for this fact was introduced in (28). Indeed, the sequence \((\langle s_t \rangle)_t\) is very similar to the sequence \((c_t)_t\) defined previously. Both move along cyclic subgroups and their definitions are very similar. However, \((\langle s_t \rangle)_t\) is not a Markov chain, whereas \((c_t)_t\) is a Markov chain. We took care of this problem by using the minimal generator \(\nu(\cdot)\) as in definition (28) of \(c_t\).

We use the Metropolis-Hastings algorithm with the above proposal distribution to approximate \(\tilde{\pi}\).

**Algorithm 15.** Starting from a permutation \(\sigma_0 \in \mathcal{G}_p\), repeat the following two steps for \(t = 1, 2, \ldots:\)

1. Sample \(x_t\) uniformly from the set \(T\) of all transpositions and set \(\sigma' = \sigma_{t-1} \circ x_t\);
2. Accept the move \(\sigma_t = \sigma'\) with probability
   \[
   \min \left\{ 1, \frac{I(\sigma') (\delta + n, D + U)}{I(\sigma) (\delta, D)} \right\}.
   \]
   If the move is rejected, set \(\sigma_t = \sigma_{t-1}\).

By the ergodicity of the Markov chain \((\sigma_t)_t\) constructed above, as the number of steps \(T \to \infty\), we have
\[
\frac{\sum_{t=1}^T 1_{\sigma = \sigma_t}}{T} \xrightarrow{a.s.} \tilde{\pi}_\sigma \quad (\sigma \in \mathcal{G}_p).
\]
This fact allows us to develop a scheme for approximating the posterior probability \(\pi\).

**Theorem 16.** We have as \(T \to \infty\),
\[
\frac{1}{\Phi(c)} \sum_{t=1}^T \frac{1}{\Phi(\sigma_t)} 1_{c = \langle \sigma_t \rangle} \xrightarrow{a.s.} \pi_c \quad (c \in \mathcal{C}).
\]

**Proof.** Let us denote \(n_{\sigma}^{(T)} = \sum_{t=1}^T 1_{\sigma = \sigma_t}, \sigma \in \mathcal{G}_p\). We have \(T = \sum_{\sigma \in \mathcal{G}_p} n_{\sigma}^{(T)}\) and \(n_{\sigma}^{(T)} / T \xrightarrow{a.s.} \tilde{\pi}_\sigma\). Moreover,
\[
\frac{\sum_{t=1}^T \frac{1}{\Phi(\sigma_t)} 1_{c = \langle \sigma_t \rangle}}{\sum_{t=1}^T \frac{1}{\Phi(\sigma_t)}} = \frac{\sum_{\sigma \in \mathcal{C}} n_{\sigma}^{(T)}}{\sum_{\gamma \in \mathcal{C}} \frac{1}{\Phi(\gamma)} 1_{\gamma = \langle \sigma_t \rangle}} = \frac{\sum_{\sigma \in \mathcal{C}} n_{\sigma}^{(T)}}{\sum_{\gamma \in \mathcal{C}} \frac{1}{\Phi(\gamma)} \sum_{\sigma \in \mathcal{C}} n_{\sigma}^{(T)}} \xrightarrow{a.s.} \frac{1}{\Phi(c)} \sum_{\sigma \in \mathcal{C}} \tilde{\pi}_\sigma \frac{1}{\sum_{\gamma \in \mathcal{C}} \frac{1}{\Phi(\gamma)} \sum_{\sigma \in \mathcal{C}} \tilde{\pi}_\gamma \pi_{\gamma},}
\]
We have \( n = 25 \) and \( p = 4 \). The head dimensions (length \( L_i \) and breadth \( B_i \), \( i = 1, 2 \)) of 25 pairs of first and second sons were measured. Thus we have \( n = 25 \) and \( p = 4 \). The following sample covariance matrix is obtained (we have \( n = 4 \)).

## Table 2

<table>
<thead>
<tr>
<th>Group ( \Gamma_k )</th>
<th>( (k_i) )</th>
<th>( (r_i) )</th>
<th>( (d_i) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( \Gamma_1 ) = {id}</td>
<td>(1)</td>
<td>(4)</td>
<td>(1)</td>
</tr>
<tr>
<td>( \Gamma_2 ) = {(1, 2)}</td>
<td>(1,1)</td>
<td>(3,1)</td>
<td>(1,1)</td>
</tr>
<tr>
<td>( \Gamma_3 ) = {(1, 3)}</td>
<td>(2,1)</td>
<td>(1,1)</td>
<td></td>
</tr>
<tr>
<td>( \Gamma_4 ) = {(1, 4)}</td>
<td>(3,1)</td>
<td>(1,1)</td>
<td></td>
</tr>
<tr>
<td>( \Gamma_5 ) = {(2, 3)}</td>
<td>(1,1)</td>
<td>(1,1)</td>
<td></td>
</tr>
<tr>
<td>( \Gamma_6 ) = {(2, 4)}</td>
<td>(1,1)</td>
<td>(1,1)</td>
<td></td>
</tr>
<tr>
<td>( \Gamma_7 ) = {(3, 4)}</td>
<td>(1,1)</td>
<td>(1,1)</td>
<td></td>
</tr>
<tr>
<td>( \Gamma_8 ) = {(1, 2, 3)}</td>
<td>(1,2)</td>
<td>(2,1)</td>
<td>(1,1)</td>
</tr>
<tr>
<td>( \Gamma_9 ) = {(1, 2, 4)}</td>
<td>(1,1,2)</td>
<td>(1,1,1)</td>
<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{10} ) = {(1, 3, 4)}</td>
<td>(1,1,2)</td>
<td>(1,1,1)</td>
<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{11} ) = {(2, 3, 4)}</td>
<td>(1,1,2)</td>
<td>(1,1,1)</td>
<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{12} ) = {(1, 2)(3, 4)}</td>
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<td>(2,2)</td>
<td>(1,1)</td>
</tr>
<tr>
<td>( \Gamma_{13} ) = {(1, 3)(2, 4)}</td>
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<td>(2,2)</td>
<td>(1,1)</td>
</tr>
<tr>
<td>( \Gamma_{14} ) = {(1, 4)(2, 3)}</td>
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<td>(2,2)</td>
<td>(1,1)</td>
</tr>
<tr>
<td>( \Gamma_{15} ) = {(1, 2, 3, 4)}</td>
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<td>(1,1,1)</td>
<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{16} ) = {(1, 2, 4, 3)}</td>
<td>(1,1,2)</td>
<td>(1,1,1)</td>
<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{17} ) = {(1, 3, 2, 4)}</td>
<td>(1,1,2)</td>
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<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{18} ) = {(1, 2, 3)}</td>
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<tr>
<td>( \Gamma_{19} ) = {(1, 3, 2)}</td>
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<td>(1,1,1)</td>
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<tr>
<td>( \Gamma_{20} ) = {(1, 4, 2)}</td>
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<td>(1,1,1)</td>
<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{21} ) = {(1, 2)(3, 4), (1, 4)(2, 3)}</td>
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<td>(1,1,1)</td>
</tr>
<tr>
<td>( \Gamma_{22} ) = (S_4)</td>
<td>(1,3)</td>
<td>(1,1)</td>
<td>(1,1)</td>
</tr>
</tbody>
</table>

Finally, by (34) we have

\[
\frac{1}{\Phi(c)} \sum_{\sigma \in C_c} \pi_\sigma = \sum_{\gamma \in C} \pi_\gamma \propto \pi_c,
\]

which completes the proof. \( \square \)

In order to approximate the posterior probability \( \pi \), we allowed the Markov chain to travel on the larger space \( S_n \). In particular, each state \( c \in C \) was multiplied \( \Phi(c) \geq 1 \) times, where \( \Phi(c) \) is the number of permutations generating \( c \). This procedure should result in slower convergence to the stationary distribution in (36). By comparing with (35), we see that (36) can be interpreted as follows: let us assign to each cyclic subgroup \( c \) a weight \( 1/\Phi(c) \leq 1 \). Then, the denominator \( N_T := \sum_{t=1}^T 1/\Phi(\langle \sigma_t \rangle) \) can be thought of as an “effective” number of steps and the numerator is the number of “effective” steps spent in state \( c \). In general, for large \( T \) we expect \( N_T \ll T \) (see an example in Section 5.2).

### 4.2. Model selection for \( p = 4 \)

Our numbering of colored models on four vertices is in accordance with (Gehrmann, 2011, Fig. 15 and 16, p. 674–675). However, we identify models by the largest group with the same coloring \( \Gamma^* \) rather than the smallest as in Gehrmann (2011). There are 30 different subgroups of \( S_4 \), which generate 22 different colored spaces. Up to conjugacy (renumbering of vertices), there are 8 different conjugacy classes. Within a conjugacy class, constants \( (k_i, r_i, d_i)_{i=1}^p \) remain the same. Groups \( \Gamma_k^* \) for \( k = 1, \ldots, 17 \) correspond to cyclic colorings.

We apply our results and methods in order to do Bayesian model selection for the celebrated example of Frets’ heads, Frets (1921), Whittaker (1990). The head dimensions (length \( L_i \) and breadth \( B_i \), \( i = 1, 2 \)) of 25 pairs of first and second sons were measured. Thus we have \( n = 25 \) and \( p = 4 \). The following sample covariance matrix is obtained (we have
MODELS INVARIANT BY SYMMETRY

1. Posterior probabilities in Frets’ heads for three best models, \( \delta = 3 \) and given \( D \).

<table>
<thead>
<tr>
<th>( D )</th>
<th>Best model</th>
<th>2nd best</th>
<th>3rd best</th>
</tr>
</thead>
<tbody>
<tr>
<td>( I_4 )</td>
<td>( \Gamma_{22} ) (95.2%)</td>
<td>( \Gamma_{16} ) (2.5%)</td>
<td>( \Gamma_{17} ) (1.3%)</td>
</tr>
<tr>
<td>( 50I_4 )</td>
<td>( \Gamma_{19}^* ) (33.8%)</td>
<td>( \Gamma_{13}^* ) (29.6%)</td>
<td>( \Gamma_{8}^* ) (13.3%)</td>
</tr>
<tr>
<td>( 100I_4 )</td>
<td>( \Gamma_{13} ) (39.6%)</td>
<td>( \Gamma_{19} ) (29.8%)</td>
<td>( \Gamma_{8} ) (7.2%)</td>
</tr>
<tr>
<td>( 1000I_4 )</td>
<td>( \Gamma_1^* ) (38.9%)</td>
<td>( \Gamma_{13} ) (10.5%)</td>
<td>( \Gamma_3 ) (10.3%)</td>
</tr>
</tbody>
</table>

\[
Z = (L_1, B_1, L_2, B_2)^	op,
\]

\[
U = \sum_{i=1}^n Z^{(i)} \cdot Z^{(i)\top} = \begin{pmatrix}
2287.04 & 1268.84 & 1671.88 & 1106.68 \\
1268.84 & 1304.64 & 1231.48 & 841.28 \\
1671.88 & 1231.48 & 2419.36 & 1356.96 \\
1106.68 & 841.28 & 1356.96 & 1080.56
\end{pmatrix}.
\]

We perform Bayesian model selection within all RCOP models, not just the ones corresponding to cyclic subgroups. In Table 2 we list all RCOP models on full graph with four vertices, along with corresponding structure constants. Structure constants remain the same within a conjugacy class, however the invariant measure \( \varphi_T \) is always different. Since there are only 22 such models, we calculate all exact posterior probabilities. The Table 2 and the invariant measures \( \varphi_T \) were obtained by using Lemma 4.

In Table 3 we summarize the results when \( \delta = 3 \) (a parameter of the prior distribution, Section 3.4), giving the three best coloring models with the highest posterior probability, for each given \( D \). Results are very similar for \( \delta = 10 \) and the given values of \( D \). For comparison, the three best models according to BIC are \( \Gamma_{19}^* \), \( \Gamma_{13}^* \) and \( \Gamma_8^* \) with the BIC 834.5, 835.4 and 835.5 respectively.

For different values of \( D = dI_4 \), the only models that have highest posterior probability are the 4 models: \( \Gamma_{22}^* = \mathcal{S}_4 \), \( \Gamma_{19}^* = \langle (1,3), (2,4) \rangle \), \( \Gamma_{13}^* = \langle (1,3)(2,4) \rangle \), \( \Gamma_1^* = \{ \text{id} \} \). These four subgroups form a path in the Hasse diagram of subgroups of \( \mathcal{S}_4^* \), i.e. \( \Gamma_{22}^* \supset \Gamma_{19}^* \supset \Gamma_{13}^* \supset \Gamma_1^* \). Thus the four selected colorings, corresponding to the permutation groups are in some way consistent. Moreover, each of them has a good statistical interpretation. Let us interpret models \( \Gamma_{13}^* \) and \( \Gamma_{19}^* \). Recall the enumeration of vertices \( (1,2,3,4) = (L_1, B_1, L_2, B_2) \). The invariance with respect to the transposition \( (1,3) \) means that \( L_1 \) is exchangeable with \( L_2 \) and, similarly, the invariance with respect to the transposition \( (2,4) \) implies exchangeability of \( B_1 \) and \( B_2 \). Both together correspond to the fact that sons should be exchangeable in some way.

We observe that only the \( \Gamma_{22}^* \) model appeared in former attempts of model selection for Frets’ heads data. It was considered in (Massam, Li and Gao, 2018, Fig. S7 p.28 of the Supplementary material) with eleven other models. Note that the only complete RCOP model selected in Gehrmann (2011) (who used the Edwards-Havránek model selection procedure) among the 9 minimally accepted models on p. 676 of her article is \( \Gamma_{10}^* \), which is not selected by our exact Bayesian procedure for any choice of \( D = dI_4 \).

5. Simulations. Let the covariance matrix \( \Sigma = (c_{ij})_{ij} \in \text{Sym}^+(p; \mathbb{R}) \) be the symmetric circulant matrix defined by

\[
c_{ij} = \begin{cases}
1 - |i - j|/p, & \text{if } i \neq j, \\
1 + 1/p, & \text{if } i = j.
\end{cases}
\]

It is easily seen that this matrix belongs to \( \mathcal{P}_{(\sigma^*\gamma)} \) with \( \sigma^* = (1, 2, \ldots, p - 1, p) \).
(a) (b) Fig 1. Heat map of matrix Σ (a) and matrix U/n (b).

Table 4

<table>
<thead>
<tr>
<th>generator of a cyclic group</th>
<th>number of visits</th>
</tr>
</thead>
<tbody>
<tr>
<td>(1, 2, 3, 4, 5, 6, 7, 8, 9, 10)</td>
<td>457725</td>
</tr>
<tr>
<td>(1, 6, 2, 7)(3, 5, 9)(4, 8, 10)</td>
<td>110677</td>
</tr>
<tr>
<td>(1, 6)(2, 7)(3, 5, 9)(4, 8, 10)</td>
<td>51168</td>
</tr>
<tr>
<td>(1, 7)(2, 6)(3, 5, 9)(4, 8, 10)</td>
<td>40895</td>
</tr>
<tr>
<td>(1, 2, 6, 7)(3, 5, 9)(4, 8, 10)</td>
<td>34883</td>
</tr>
</tbody>
</table>

Table 5

<table>
<thead>
<tr>
<th>generator of a cyclic group</th>
<th>#most visited</th>
<th>ARI</th>
</tr>
</thead>
<tbody>
<tr>
<td>(1, 2, 3, 4, 5, 6, 7, 8, 9, 10)</td>
<td>25</td>
<td>1.00</td>
</tr>
<tr>
<td>(1, 3, 5, 7, 9)(2, 4, 6, 8, 10)</td>
<td>13</td>
<td>0.60</td>
</tr>
<tr>
<td>(1, 2, 4, 3, 5, 6, 7, 9, 8, 10)</td>
<td>3</td>
<td>0.43</td>
</tr>
<tr>
<td>(1, 2, 4, 3, 5, 6, 7, 9, 8, 10)</td>
<td>2</td>
<td>0.46</td>
</tr>
<tr>
<td>(1, 3, 2, 4, 5, 6, 8, 7, 9, 10)</td>
<td>2</td>
<td>0.43</td>
</tr>
<tr>
<td>(1, 3, 5, 9, 2, 6, 8, 10, 4, 7)</td>
<td>2</td>
<td>0.43</td>
</tr>
<tr>
<td>(1, 2, 4, 3, 5, 6, 9, 8, 10, 7)</td>
<td>2</td>
<td>0.35</td>
</tr>
<tr>
<td>(1, 4, 5, 7, 8)(2, 3, 6, 9, 10)</td>
<td>2</td>
<td>0.24</td>
</tr>
<tr>
<td>(1, 8, 10, 9)(2, 7)(3, 5, 4, 6)</td>
<td>2</td>
<td>0.19</td>
</tr>
<tr>
<td>(1, 2, 10, 3)(4, 9)(5, 8, 6, 7)</td>
<td>2</td>
<td>0.19</td>
</tr>
</tbody>
</table>

5.1. First approach. For \( p = 10 \) and \( n = 20 \), we sampled \( Z^{(1)}, \ldots, Z^{(n)} \) from the \( N_p(0, \Sigma) \) distribution and obtained \( U = \sum_{i=1}^{n} Z^{(i)} \cdot Z^{(i)^T} \) depicted in Fig. 1 (b).

We run the Metropolis-Hastings algorithm starting from the group \( \langle \sigma_0 \rangle = \{\text{id}\} \) with hyper-parameters \( \delta = 3 \) and \( D = I_{10} \). After 1 000 000 steps, the five most visited states are given in the Tab. 4.

The Metropolis-Hastings (M-H) algorithm recovered the true pattern of the covariance matrix. The acceptance rate was 2.5\% and the Markov chain visited 746 different cyclic groups. The acceptance rate can be increased by a suitable choice of the hyper-parameters (e.g. for \( D = 10I_{10} \) the acceptance rate is around 10\%).

In order to grasp how randomness may influence results, we performed 100 simulations, where each time we sample \( Z^{(1)}, \ldots, Z^{(n)} \) from \( N_p(0, \Sigma) \) and we run M-H for 100 000 steps with the same parameters as before. In Table 5 we present how many times a given cyclic subgroup was most visited during these 100 simulations (second column). There were 53 distinct cyclic subgroups, which were most visited at least in one of the 100 simulations; below we present 10 such subgroups. The average acceptance rate is 1.4\% (see the histogram in Fig. 2). When we regard colorings as partitions of the set \( V \cup E \) according to group orbit decomposition, the two colorings may be compared using the so-called adjusted Rand index.
The order of \( c \) equal to (recall (36)), a similarity measure comparing partitions which takes values between \(-1\) and 1, where 1 stands for perfect match and independent random labelings have score close to 0. In the third column of Table 5, we give the adjusted Rand index between the colorings generated by given cyclic subgroup and the true coloring.

We see that groups which were most visited by the Markov chain have positive ARI and the true pattern was recovered in a quarter of cases. We stress that even though the colorings generated by \((1, 2, 3, 4, 5, 6, 7, 8, 9, 10)\) and \((1, 3, 5, 7, 9)(2, 4, 6, 8, 10)\) are very similar, the distance between these subgroups is 9, that is, the Markov chain \((C_t)_t\) needs at least 9 steps to get from one subgroup to the other. We performed similar simulations for \(n = p = 10\) and the results were only slightly worse: the true pattern was recovered in 18 out of 100 runs of the algorithm.

This indicates that the Markov chain may encounter many local maxima and one should always tune the hyper parameters in order to have higher acceptance rate or to allow the Markov chain \((C_t)_t\) to make bigger steps.

5.2. Second approach. We performed \(T = 100000\) steps of Algorithm 15 with \(\sigma_0 = \text{id}, p = 100, n = 200, \delta = 3\) and \(D = I_{100}\). Let us note that for \(p = 100\), there are about \(4 \cdot 10^{155}\) cyclic subgroups and this is the number of models we consider in our model search.

We have used Theorem 16 to approximate the posterior probability distribution \((\pi_c; c \in C)\) (see (30)). The highest estimated posterior probability was obtained for \(c^* := \langle \sigma^* \rangle\), where

\[
\sigma^* = (1, 2, 3, 4)(6, 8, 15)(7, 10, 9)(11, 16, 12)(13, 17, 14)(18, 19, 20, 22, 21)(23, 26)
\]

\[
\]

\[
(38, 47, 41, 48)(43, 51, 46, 49)(50, 52, 53, 54)(56, 58, 57)(59, 66, 67)
\]

\[
(60, 65, 63)(61, 62, 64)(68, 71, 72, 70, 69)(73, 93)(74, 77)(75, 98, 81, 100)
\]

\[
\]

The order of \( c^* \) is \(|c^*| = 60\) and \(\Phi(c^*) = 16\). The estimate of the posterior probability \(\pi_{c^*}\) is equal to (recall (36))

\[
\frac{1}{\Phi(c^*)} \sum_{t=1}^{T} \frac{1}{\Phi(\sigma_t)} \approx \frac{2361.5}{6381.5} \approx 37\%.
\]

The true covariance matrix \(\Sigma\), the data matrix \(U/n\) and the projection \(\Pi_{c^*}(U/n)\) are illustrated in Fig. 3.

We visualize the performance of the algorithm on Fig 4. In red color, a sequence \(\left(\sum_{t=1}^{k} \frac{1}{\Phi(\sigma_t)}\right)_k\) is depicted, which can be thought of as an “effective” number of steps of the
algorithm (for an explanation, see the paragraph at the end of Subsection 4.1.2). In blue, we present a sequence \( \left( \sum_{t=1}^{k} \frac{1}{\Phi(\langle \sigma_t \rangle)} \mathbf{1} (\langle \sigma_t \rangle \neq \langle \sigma_{t-1} \rangle) \right)_k \), which represents the number of weighted accepted steps, where the weight of the \( k \)th step equals \( \frac{1}{\Phi(\langle \sigma_k \rangle)} \). We restricted the plot to steps \( k = 1, \ldots, 10000 \), because after 10000 steps, the Markov chain \( (\sigma_t)_{10000 \leq t \leq 100000} \) changed its state only 9 times. For \( k = 100000 \), the value of the blue curve is 25.75, while the value of red one is 6381.5.

The model suffers from poor acceptance rate, which could be improved by an appropriate choice of the hyper-parameter \( D \) or by allowing the Markov chain to do bigger steps.

**Acknowledgments.** The authors would like to thank Steffen Lauritzen for his interest and encouragements. We also thank M. Bogdan from Wroclaw University, A. Descatha from INSERM and Centre Hospitalier Universitaire Angers and V. Seegers from Institut de Cancérologie de l’Ouest Nantes for explaining the specific nature of medical and genetic data. The paper benefited from the comments of an anonymous referee to whom the authors are grateful.

**SUPPLEMENTARY MATERIAL**

Supplement to “Model selection in the space of Gaussian models invariant by symmetry”.
Supplement contains proofs and examples. We provide proofs of Theorems 1, 5, 6 along with a background on representation theory that is needed to understand proofs. Moreover, we present proofs of Proposition 7, Theorems 8 and 9, an example to Section 2.3, proof of Lemma 13 and the real data example considered in Miller et al. (2005) and Højsgaard and Lauritzen (2008).
REFERENCES


GRACZYK, P., ISHI, H., KOLODZIEJEK, B. and MASSAM, H. Supplement to “Model selection in the space of Gaussian models invariant by symmetry”.


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