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# PUBLIC WASTE COLLECTION: A CASE STUDY 

## Ludo F. GELDERS <br> and <br> Dirk G. CATTRYSSE

Katholieke Universiteit Leuven
Afdeling Industrieel Beleid Celestijnenlaan 300 A, B-3001 Leuven-Heverlee, Belgium

## ABSTRACT

This paper deals with the waste collection in the N.E. area of Brussels. It summarizes a study carried out as a master's thesis project. Until now the collection scheme was based upon experience. Management felt that a more normative and systematic approach was needed. This paper discusses the modeling of the real-life problem based on the capacitated arc routing problem. The solution approach is based on the path scanning algorithm. The problem solver was coded in Pascal and linked with dBase-files which contain all information on the collection area. A reduction of approximately $15 \%$ in distance travelled was achieved.

## 1. PROBLEM DESCRIETION.

The current problem deals with the waste collection of a set of 5 municipalities in the N.E. area of Brussels. Between 4 and 11 waste collection trucks have to serve these communities daily. On a yearly basis, these lorries total 360.000 km . These 360.000 km are divided into 60.000 km for the actual waste collection and 300.000 km for trips to the waste disposal area, which is located south of Brussels. The distance from the last collection point to the disposal area varies between 25 and 50 km .

Until now, the collection scheme was based upon experience. Each truck was assigned a set of streets to be served. In the current system, a truck crew is free to go home whenever the job is finished. This work organization relies upon the hypothesis that the work force will try to optimize its routing. Under current operations, the company employs 41 people, 34 being crew members of the trucks.

As always in garbage collection problems, the experience based work schemes take into account a very broad and complicated set of constraints, e.g. :

- one-way streets
- fluctuating quantities to be collected
- given collection frequency (twice per week)
- capacity limitation of trucks
- truck maintenance
- special collection services due to special events
- etc...etc...

Management felt however that a more normative and systematic approach might give some insight to do the job more efficient, through a better use of resources (decreasing the total travel distance, the number of trucks used, etc...). Moreover a need was felt to have an interactive decision tool which might provide quick answers to certain urgent problems (e.g. unavailability of a crew or a truck).

## 2. MODELING THE PROBLEM.

In literature one can find discussions on similar problems, e.g. : sanitation vehicle routing [ 2], vehicle routing for municipal waste collection [ 1] and routing electric meter readers [ 9].

The above waste collection problem can be modelled as the Capacitated Arc Routing Problem (CARP) : given an undirected network $G(N, E, C)$ with arc demands $q_{i j} \geq 0$ for each arc (i,j) which must be satisfied by one of a fleet of vehicles of capacity $W$, find a number of cycles each of which passes through the depot (node 1) which satisfy demands at minimal total cost. The CARP can be formulated as follows (see Golden and Wong [ 7]) :
$\min \sum_{i=1}^{N} \sum_{j=1}^{N} \sum_{k=1}^{K} \quad c_{i j} \quad x_{i j k}$

$$
\text { s.t. } \begin{array}{ll}
\sum_{j=1}^{N} x_{j i k}-\sum_{j-1}^{N} x_{i j k}=0 & \begin{array}{l}
i=1, \ldots, N \\
k
\end{array}=1, \ldots, K
\end{array}
$$

$$
x_{i j k} \geq 1_{1 j k}
$$

$$
\sum_{i=1}^{N} \sum_{j=1}^{N} 1_{i j k} q_{i j} \leq W \quad k=1, \ldots, K
$$

$$
\sum_{i \in Q} \sum_{j \in Q} x_{i j k}-N^{2} y_{l q k} \leq|Q|-1
$$

$$
\sum_{i \in Q} \sum_{j \notin Q} x_{i} j k+y_{2 q k} \geq 1
$$

$$
y_{1 q k}+y_{2 q k} \leq 1
$$

$$
y_{1 q k}, y_{2 q k}, x_{i j k}, 1_{i j k} \in\{0,1\}
$$

```
k = 1, ., K
```

$\mathrm{q}=1, \ldots, 2^{\mathrm{N}-1}-1$
and every nonempty
subset $Q$ of $\{2,3, \ldots, N\}$
where $N=$ the number of nodes,
$K=$ the number of available vehicles,
$q_{i j}=$ the demand on arc (i,j),
$W=$ the vehicle capacity $\left(W \geq \max q_{i j}\right)$,
$c_{i j}=$ the length of arc (i,j),
$x_{i j k}=1$, if arc (i,j) is traversed by vehicle $k$, 0 otherwise,
$1_{i j k}=1$, if vehicle $k$ services arc (i,j), 0 otherwise,
$\lceil z\rceil=$ the smallest integer greater than or equal to $z$,
$Q=$ subset of $N$ with cardinality between 2 and $|N|$.

The objective function (1) seeks to minimize total distance travelled. Equations (2) ensure route continuity. Equations (3) state that each arc with positive demand is serviced exactly once. Equations (4) guarantee that arc ( $i, j$ ) can be serviced by vehicle $k$ only if it covers arc (i,j). Vehicle capacity is not violated on account of equations (5). Equations (6) prohibit the formation of (infeasible) subtours. Every index q corresponds to a set $Q$. Integrality restrictions are given in (7).

## 3. SELECTION OF SOLUTION METHODOLOGY.

Given the computational complexity of the CARP (which is an NP-hard problem (Golden and Wong [ 7])) it becomes necessary to apply approximate solution techniques or heuristics. Golden, De Armon and Baker [ 6] compare some heuristics developed to solve the CCPP (Capacitated Chinese Postman Problem). The CCPP is a special form of the CARP : $q_{i j}>0$ for all $i$ and $j$ (instead of $q_{i j} \geq 0$ ). In general, 3 different approaches are available when dealing with multi-vehicle problems , i.e.,

- Cluster First-Route Second : the area is divided in subareas in such a way that every subarea can be serviced by one vehicle. For every subarea a Chinese Postman Problem (CPP) is solved.
- Route First-Cluster Second : first the problem is solved as a CPP. The complete cycle is then split in several routes which do not violate the capacity constraints of the vehicles and minimize the distance travelled.
- Route and Cluster Together : routing and clustering is done at the same time.

All three approaches have advantages and disadvantages. The first method creates smaller problems which are easy to solve and requires less data handling but yields suboptimal solutions and does not take the stochastic character of the demand into account. The second method has the advantage that the first phase remains the same even if demand varies. However a larger network has to be optimized and the relation between the touis and the depot (start node) and the different tours is gone (remember that the distance to the depot contributes the most to the total travelling distance). This method works rather well when every tour consists of only a few arcs [ 9]. This is not the case in this study. The third approach is well adapted to problems where a vast physical area contains a very large number of arcs. It allows for a good overall and integrated view of the problem. This third approach is chosen and three algorithms in that class are briefly discussed :

- Path Scanning [ 6] : the basic idea is to construct a tour at a time by adding arcs sequentially till the capacity is exhausted. Then the shortest return path to the depot is followed. The criteria for choosing an arc (i,j) are : the distance, $c_{i, 1}$, per unit remaining demand is minimized (a) or maximized (b), the distance from node $j$ back to the depot is minimized (c) or maximized (d), if the vehicle is less than half-full maximize the distance (e); (a) looks for a large and quick payoff while (b) seeks to incur the larger expenses early; (c) tends to obtain shorter cycles; whereas (d) in general, yields longer cycles; (e) represents a hybrid approach. The set of cycles with smallest total distance is selected as outcome for this simple and rather fast algorithm.
- Construct and Strike [ 4, 6] : this algorithm repeatedly constructs feasible cycles and then strikes or removes them. The following steps are required :
step 1 : construct a cycle fulfilling the capacity constraint of the vehicle.
step 2 : set the demand of the arcs belonging to this cycle equal to 0 .
step 3 : repeat steps 1 and 2 till there are no arcs, to be serviced, left in the network.

This procedure has one drawback : how to construct good cycles ?

- Augment - Merge [ 3, 7] : this procedure merges several smaller cycles
in larges ones :
step 1 : start with as much cycles as arcs to be served.
step 2 : starting with the largest cycle available, see if a demand arc
on a smaller cycle can be serviced on a larger cycle
( $=$ AUGMENT).
step 3 : evaluate the merging of any two cycles, subject to capacity
constraints or additional restrictions. Merge the two cycles
which yield the largest positive savings ( $=$ MERGE).
step 4 : repeat step 3 until finished.
Computational experience done by Golden, De Armon and Baker [ 6]
indicates that the augment-merge strategy yields the best results on
average, but is much more complicated and time consuming than the two
other types. The path scanning algorithm has the advantages of
simplicity and minor CPU-time requirements. This algorithm can easily be
adapted to all the necessary needs and criteria as required by the
problem stated above. The construct and strike algorithm has the same
advantages but yields results which are worse.
Since management needs an interactive decision tool which provides quick
answers to urgent problems, an approach based on the path scanning
algorithm is chosen. Moreover, a number of different criteria for
choosing the next arc in every step of the algorithm exist (they are
mentioned earlier). By simulating some scenarios it is found that the
following combination of two criteria leads to the best results.
Given a provisional path (start node $=1$ ) arrives in node i, we add the
arc (i,j) satisfying one of the following criteria :
- if the truck load is less than $50 \%$, maximize the shortest distance
from $j$ to the end-node, distance of arc ( $i, j$ ) included.
- if the truck load is $50 \%$ or more, minimize the shortest distance from
$j$ to the end-node, distance of arc (i,j) included.
This combination leads to cycles with a reasonable length without
augmenting markedly the distance from the last collection point back to
the start or end (waste disposal) node.


## 4. IMPLEMENTATION ISSUES.

### 4.1 CLUSTERING.

Although the path scanning algorithm is a route and cluster together method, we will start dividing the total collection area in a number of sectors. These sectors are much larger than the clusters in the cluster first - route second approach. For a particular day, all lorries will be send to the same sector. Our major concern is to minimize the number of trips from and to the waste disposal node. In this way a lorry can take over a part of the tour of another one without driving a long distance (interaction between tours). Demand is stochastic and a vehicle can reach its maximum allowable load before finishing its cycle. Adding up the collection amounts to obtain the total daily collection quantity per sector, levels out a large number of these stochastic elements. Nevertheless some time-depending parameters remain. There is a seasonal dependency because people produce more garbage in summer time than in winter. Also on the long term, the total amount of garbage to collect yearly increases. To start, we will define the sectors by aggregating the present sectors.

### 4.2 INPUT DATA.

It was an enormous task to collect all input data necessary for solving this problem. And as the accuracy of these data will influence to a large extent the accuracy of the solution, it is of major concern to keep these inputs up-to-date. So a well developed database structure had to be set up. The input data can be devided into three classes:

- data related to the network :
- an adjacency list for every node
- cost $c_{1 j}$ for every arc, expressed as:

1. a collection cost (time)
2. a driving cost (time to drive through the arc without collection)

- quantity garbage $q_{i j}$ to be collected for arc (i,j)
- one-way streets or not
- both sides of the streets can or can not be served simultaneously
- collection arcs and arcs connecting one node to another without collecting garbage
- start and end node for every sector.
- number of lorries available for the different types of garbage
- maximum capacity (load) for every lorry and an average lorry capacity $W$ as used in the program
- as the density of garbage is not always the same, not weight but volume can be the binding constraint. We do not only need an average $W$ but also a margin on $W$, to express the variability
- seasonal coefficient on quantity of garbage produced.
- a number of parameters tested in the different scenarios to evaluate their sensitivity on the result.

The most important data in this model are the costs $c_{i j}$ and the quantities $q_{i j}$ related to every arc. $c_{i j}$ can be expressed as the time necessary to serve an arc, with or without collection. The time can be derived by measuring the distance $\left(l_{1 j}\right)$ and multiplying it by the average velocity of a lorry. Distances of the arc were measured on detailed maps of the area and corrected with information from the drivers who noted down distance travelled. There are 2 velocities:

$$
\begin{array}{lll}
v_{d}=\text { velocity with collection } & --> & t_{d i f}=l_{i j} / v_{d} \\
v_{b}=\text { velocity without collection } & --> & t_{b i j}=l_{i j} / v_{b}
\end{array}
$$

In the article of Bodin et al. [ 2] the velocity $v_{b}$ is even devided into 4 classes depending on the types of streets.
The quantity $q_{i j}$ is the most difficult element in gathering the data, because $q_{i j}$ has a stochastic nature and a direct measurement (estimation) of $q_{i j}$ is impossible for practical reasons.
$q_{1 f}$ was estimated in the following way. From historical data, total weakly collection volumes $V$ per sector $S$ will be determined.
$q_{i j}=f_{i j} * V$ with $0 \leq f_{i j} \leq 1$,
( $f_{1 j}=0$ if arc ( $i, j$ ) belongs to the sector $S$ and has to be serviced by another sector $S^{\prime}$ or arc (i,j) does not belong to the sector $S$ )
When garbage is collected in mini-containers, $f_{i j}$ can be easily estimated as follows :
$f_{i j}=\frac{\text { \# containers for } \operatorname{arc}(i, j)}{\# \text { containers for sector } S(\operatorname{arc}(i, j) \in S)}$

It is clear that the costs $c_{i j}$ can be estimated more accurately than the collection quantities $q_{i j}$ per arc. Nevertheless, with regard to the accuracy of the data, following equation is very important to keep in mind :

$$
\begin{aligned}
e_{r}(q)=e_{r}\left(q_{i j}\right) * n^{-1 / 2} \text { with: } q & =\text { collection quantity per trip } \\
& n=\text { number of arcs in one trip } \\
& e_{r}(X)
\end{aligned}
$$

Analogously we have:
$e_{r}(1)=e_{r}\left(1_{i j}\right) * n^{-1 / 2} \quad$ with: $\quad=$ driving distance per trip
$1_{i j}=$ driving distance per arc

Example: $\quad 1=20.000 \mathrm{~m}$
$1_{i j}=200 m \quad-\ldots->e_{r}(1)=e_{r}\left(1_{i j}\right) * 10$
$\mathrm{n} \quad=1 / 1_{\mathrm{ij}}=100$

This means that when a relative error of $5 \%$ on 1 with a confidence interval of $95 \%$ is required, a relative error on $l_{i j}$ up to $50 \%$ is allowed. As a consequence the way we estimated both distances and collection quantities per arc leads to rather accurate results when aggregated over a collection trip.

## 5. COMPUTATIONAL ASPECTS.

The software program can be separated into two parts:

- a general database program (in dBase III), to store and update al1 input data.
- a route generation program written in Pascal.
5.1 GENERAL DATABASE PROGRAM.

Every arc is defined by two nodes, the start- and end-node. The following information is given per arc :

- two costs : a collection cost and an empty cost (cost for driving through an arc without collecting garbage)
- capacity $q_{i j}$ : quantity of garbage to collect for arc (i,j)
- type : $1=$ directed arc
$2=$ undirected arc
- ten fields $S 1$, $S 2 \ldots . S 10$, corresponding to ten sectors with each field containing one of the following codes : $1=a \operatorname{collection~arc~for~the~corresponding~sector.~}$ $2=$ arc belongs to the sector, but no collection is necessary. $3=$ arc does not belong to the sector.
This is a general database, containing all the possible streets in the collection problem. From this data base we can easily select the arcs belonging to one or more sectors (1 ... 10), these arcs will be copied into a sector data base and used as input for the route generation program. Whenever there is a change in the street pattern or collection quantities, the general data base will be updated and new partial data bases per sector will be derived. Organizing the input structure in this way will never lead to data inconsistency.


### 5.2 THE ROUTE GENERATION PROGRAM.

The program is built in a modular way. The first module reads all the network data out of the sector data base and doubles the undirected arcs. The second module reads the general data related to the sector (startnode, end-node, collection quantity) and other data like number of lorries, their capacity, etc...

After the user has chosen an optimization criterion, one arc at a time will be selected in module 3 to generate the different trips. In every node the path scanning algorithm is used when one of the related arcs is not yet served. The Dijkstra [ 5] algorithm is used in all other cases, to find the shortest way to another unserved arc.
When lorry capacity is reached, the program also uses the Dijkstra algorithm to find the shortest way back to the end-point.
A new trip will start in the end- or start-node of the network and module 3 will again generate a new routing. New trips will be generated as long as unserved arcs in the network exist. After every trip, data with regard to collection quantity $q$ and length of the trip (l) are stored.

### 5.3 FLEXIBILITY OF THE PROGRAM.

The program is built in such a way that it is very easy to generate a number of scenarios. Following options are possible :

- changing the area and shape of the sectors
- choosing a start- and end-point of a trip in one node or in two different nodes. In this way a lorry can start its first trip at the depot and end it at the waste collection area (from where the second trip starts)
- different criteria to execute the path scanning algorithm
- introduction of a seasonal coefficient to adapt collection quantities per arc over the year
- a lorry, close to the end-point of the network has two options. It can start a new cycle, or it can drive directly to the waste collection area. A preference for one of these options can be implemented by adding a parameter $\alpha$. Instead of having a lorry capacity $W$, we will give the truck in the end-point a capacity $\alpha W$. If $q<\alpha W$, the truck will start a new cycle. If $q>\alpha W$, the truck will go directly to the collection area.

By giving $\alpha$ a number between 0 and 1 we can influence thoroughly the length and the number of trips

- simultaneous garbage collection on both sides of the street
- blind alleys are automatically added to an adjacent arc. In this way the network is simplified and empty driving time is minimized.

Elements not included in the program are : time constraints, holidays, rough garbage, traffic jams, right turns.

## 6. CONCLUSIONS

A combination of six current sectors was selected as a testing area.
These sectors give a good representation of the whole collection area. Total weekly collection in the testing area is approximately 5000 tons.

## 6. 1 INFLUENCE OF LORRY CAPACITY (W).

Initially $W$ was fixed on 9.330 Kg . The total driving distance decreases by increasing $W$. From fig. 1 we see that the decrease in driving distance is not continuous. The explanation is very simple. Whenever a lorry A reaches its capacity constraint $W$, it drives directly to the waste disposal area. Another lorry B (operating in the neighbourhood) will continue and serve the remaining collection quantity of the arc (i,j). Increasing $W$ will only decrease total driving distance when arc (i,j) can be collected completely by lorry A.


Fig.l : Sensitivity of the lorry capacity on the total driving distance.

### 6.2 OPTIMAL CRITERION FOR OPEN AND CLOSED TRIPS.

A closed trip starts and ends in the same point of the network whereas, an open trip has a different start- and end-point. Intuitively it is clear that open trips give better results than closed trips for all possible criteria. However for combination tours, open as well as closed trips, using one and the same criterion for both types of trips gives better results than using two different criteria. This was a rather remarkable result we found from the simulations of the testing area.

### 6.3 INFLUENCE OF PARAMETER $\alpha$.

As mentioned before, parameter $\alpha$ is of major influence on the number of trips and on the collection distance. The smaller $\alpha$, the larger the number of trips (trips are shorter). It is more important to reduce the number of trips than to reduce the total collection distance, as the waste collection point is situated a long way from the collection area. Therefore parameter $\alpha$ close to 1 yields the best results in this case.

### 6.4 COMPARISON WITH THE CURRENT SYSTEM.

Using the aggregated collection area of the six current sectors, an overall reduction in driving distance of $15 \%$ was found. Due to a
reduction in the number of trips (10 trips/week against minimum 12 trips/week in the current system) total driving distance to and from the collection area was reduced by $14.5 \%$. Another $4.5 \%$ was saved on collection distance. For the individual sectors there was no clear improvement. Whereas the Path Scanning Algorithm was developed for a CCPP problem, the scenario of individual sectors is more related to a CPP problem as one trip is mostly sufficient for collecting all waste in a sector. Aggregating sectors gives a better overall result as it reduces the number of partly loaded trucks which has an influence on the total number of trips.

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# THE STATIONARITY OF CAPM-BETA IN A CHANGING ECONOMIC ENVIRONMENT 

Winfried G. HALLERBACH
Department of Finance
Erasmus University Rotterdam
P.O. Box 1738, NL-3000 DR Rotterdam, The Netherlands

## ABSTRACT

In this paper, we specify a theoretical multi-factor return generating process of securities by means of which fundamental shifts in the structure of the economy can be modelled. We then analyze the conditions under which CAPM's risk measure 'beta' will be constant in this changing economy. These conditions are very restrictive and not likely to be met in practice.

## 1. INTRODUCTION

The concept of systematic risk is of central importance in Modern Portfolio Theory (MPT). Within the standard Capital Asset Pricing Mode1 (Sharpe [1964], Lintner [1965], Mossin [1966]), the measure for this market risk is defined as:

$$
\begin{equation*}
\beta_{i}=\frac{\operatorname{Cov}\left(r_{1}, r_{m}\right)}{\operatorname{Var}\left(r_{m}\right)} \tag{1}
\end{equation*}
$$

where $r_{1}=$ the return on security $i$;
$r_{m}=$ the return on the market portfolio $m$.
Cov and Var denote the covariance and variance operator, respectively.
In applications of MPT, historical data are used to estimate the relevant ex ante $\beta$ by means of the market model (cf. Fama [1976]). In its most simple form, this regression model can be written as:

$$
\begin{equation*}
r_{i t}=\alpha_{i}+\beta_{i} r_{m t}+\varepsilon_{i t}, \quad t \in T, i \in N \tag{2}
\end{equation*}
$$

where $\alpha_{i}$ is a constant and $\varepsilon_{i t}$ is a zero-mean error term, uncorrelated with $r_{m t}$. $T$ and $N$ denote the sample of time periods and the set of $N$ securities in the market portfolio, respectively. For the sake of brevity, the time index $t$ will be suppressed.

The use of this market model involves a number of more or less arbitrary decisions, for example the choices with respect to:
-the particular form of the model: the simple form eq. (2), the excess return (or risk premium) form or the two factor model (cf. Blume \& Friend [1973]);

- the index that serves as proxy for the market portfolio (for the effect of index choice, see for example Frankfurter [1976] and Rol1 [1977]); -the total interval $T$ chosen to perform the regression (cf. Baesel [1974], Levhari \& Levy [1977] and Theobald [1981]); the appropriate length of the observation intervals within the sample of time periods $T$ (cf. Hawawini [1983] and Handa, Kothari \& Wasley
[1989]).
The sensitivity of the estimated $\beta$ for the latter intervalling aspect determines its stability. Especially when daily data are used, nonsynchronous trading gives rise to an errors-in-variables problem and causes biased estimators (Scholes \& Williams [1977], Smith [1978], Scott \& Brown [1980], Dimson \& Marsh [1983] and more recently Shanken [1987]).

In addition, as the ex ante value of $\beta$ is relevant, it is important whether $\beta$ is constant over time. This brings us to the issue of the intertemporal stationarity ${ }^{1)}$ of $\beta$ : is a $\beta$ estimated from data in period $T_{1}$ equal to the $\beta$ estimated in period $T_{2}$ ?

Blume [1971, 1975] found that $\beta^{\prime}$ s were non-stationary. More specific, he found that, over time, $\beta^{\prime}$ s tend to drift to the market average of one. Klemkosky \& Martin [1975] examined various techniques by means of which $\beta^{\prime}$ s could be adjusted for this drift. Elgers et al. [1979] however, showed that whether $\beta^{\prime}$ s were calculated moving forward or backward in time, they similarly drifted towards the value one. This $\beta$-drift thus appears to be caused by statistical aberrations.

Fabozzi \& Clark [1978], Sunder [1980], Chen [1981] and Simonds et al. [1986], among others, test a random coefficient model and also present evidence of non-stationarity. Fisher \& Kamin [1985] suggest some improvements in estimating $\beta^{\prime}$ s by gradually decreasing the weight of each observation with the passage of time. To determine these weights, they use a Kalman filter.

Other studies only focus on part of the $\beta$-coefficient in eq. (1). Myers [1973] finds significant changes in the variability of the 'market factor' from period to period and concludes that, as a result, $\beta^{\prime} s$ are not stationary. Elton et al. [1978] focus on the numerator and investigate the stationarity of the correlation structure of security returns. They conclude that the best way to estimate securities' correlations is just to use the average correlation coefficient for the entire universe of securities ('overall mean'): additional efforts to refine the estimated correlations appeared to be fruitless.

1) Note that we explicitely make a difference between stability and stationarity.

The above mentioned studies are of an empirical nature and concentrate on the form and statistical properties of the instationarity. On the other hand, there are studies of a theoretical nature that concentrate on the identification of micro- or macroeconomic variables that influence $\beta$-stationarity. Officer [1973] investigated changes in the variability of the market factor (just as Myers [1973] did) and found that these changes appear to be related to changes in the Index of Industrial Production (a measure of business activity in the US) and changes in the M2 money supply. Levy [1971] related shifts in $\beta$ to alternating bull and bear market conditions. Fabozzi \& Francis [1977] reject the influence of bull-bear market forces on $\beta$, but conclude later that the intertemporal non-stationarity of $\beta$ appears to result from changes which are associated with business cycle economics (Francis \& Fabozzi [1979]). Using the concept of duration, Bildersee \& Roberts [1981] relate changes in $\beta$ to changing interest rates. More recently, DeJong \& Collins [1985] start from the joint Option Pricing Model/ Capital Asset Pricing Model framework and find a significant relation between $\beta$-non-stationarity and the degree of leverage of the respective firm and changes in the risk-free interest rate.

McDonald [1985] considers major structural changes in the economy. These changes were accompanied by large changes in interest rates and yielded shifts in $\beta^{\prime}$ s. Further evidence concerning these changes in the structure of the economy and the variance-covariance structure of security returns is presented by Bollerslev et al. [1988], Harvey [1989] and Van Der Meulen [1987, 1989]. Analyzing the covariance structure of deflated returns of general classes of assets (including currencies), Van Der Meulen detects several large shifts in the covariances between these assets over time. As these changes affect the composition of the assets' covariance matrix, the $\beta^{\prime}$ s of these asset classes (and of the individual assets therein) presumably will not be constant.

In this paper, we accept the strong empirical evidence supporting the existence of $\beta$-non-stationarity, and consider the theoretical effect of fundamental shifts in the structure of the economy on $\beta$. We hypothesize that the returns of assets or securities are generated by a
factor model and that intertemporal changes in the variances of the factors provide the 'missing link' between the shifts in the covariance structure of these assets. In section two, we specify the theoretical multi-factor return generating process of securities by means of which the fundamental changes in the structure of the economy can be modeled. The conditions under which security- $\boldsymbol{\beta}^{\prime} \mathrm{s}$ will be constant in a changing economy form the central issue of this paper and will be dealt with in section three. Section four contains our conclusions and directions for future research

## 2. hULTI-FACTOR MODELS AS RETURN GENERATING PROCESSES

To model the dynamic ecomomic environment in which the returns of the securities are generated, we reduce the dimensionality of the interrelationships therein. We therefore introduce a mechanism that generates the return for each of the N securities in the market portfolio. More specific, we assume that there exists a function $g_{i}($. that links the individual returns $r_{i}$ to a set of $K$ common underlying variables or factors $\left\{\delta_{j}\right\}_{j \in K}$. In a general form, this return generating process (henceforth: RGP) can be expressed as:

$$
\begin{equation*}
r_{i t}=g_{i}\left(\delta_{j t}\right)_{j \in K}+\epsilon_{i t}, \text { with } E\left(\epsilon_{i t}\right)=0, i \in N \tag{3}
\end{equation*}
$$

To incorporate randomness, the common factors are assumed to be stochastic and a random disturbance term $\epsilon_{i}$ is added to incorporate the influence of idiosyncratic factors on each individual return. We assume that the idiosyncratic factors $\epsilon_{i}$ and the common factors $\left\{\delta_{j}\right\}_{j \in K}$ are mutually independent. Furthermore, the idiosyncratic factors are assumed to be truly security specific and thus also mutually independent. With these assumptions, the dimensionality of the economic environment is reduced from ${ }^{1} \mathrm{~N}(\mathrm{~N}-1)$ relationships between the returns of securities to (i) the ${ }_{2} \mathrm{~K}(\mathrm{~K}-1)$ relationships between the common factors and (ii) the N $\times \mathrm{K}$ relations between the securities and these factors. The latter relations can be explicitized by means of sensitivity coefficients.

Although the function $g_{1}($.$) is unknown, we can apply the multi-$ variate version of Taylor's theorem and expand $g_{i}($.$) , for example,$ around the value it takes when the common factors are set to their mean values. We then can rewrite eq. (3) as

$$
\begin{align*}
r_{i}=g_{i}\left[E\left(\delta_{j}\right)\right] & +\sum_{j}\left[\delta_{j}-E\left(\delta_{j}\right)\right] \frac{\partial g_{i}}{\partial \delta_{j}} \\
& +\sum_{h=2}^{\infty} \frac{1}{h!}\left[\sum_{j}\left[\delta_{j}-E\left(\delta_{j}\right)\right] \frac{\partial}{\partial \delta_{j}}\right]^{h} g_{j}+\epsilon_{i} \tag{4}
\end{align*}
$$

where the symbolic power between the large square brackets on the right is first to be expanded formally by the binomial theorem and then the powers $\partial / \partial \delta_{j}$ multiplied by $g_{i}$ are to be replaced by the corresponding $n^{\text {th }}$ derivatives $\partial^{n} g_{i} / \partial \delta_{1}{ }^{n}, \partial^{n} g_{i} /\left(\partial \delta_{1}^{n-1} \partial \delta_{2}\right) \& c .$, evaluated at the spanning point $\left[E\left(\delta_{1}\right), E\left(\delta_{2}\right), \ldots\right]$.

Note that $g_{i}\left[E\left(\delta_{j}\right)\right]$ is not necessarily equal to $E\left(r_{i}\right)$ when $g_{i}$ is a non-1inear function of the $\delta_{j}$ 's (cf. Jensen's inequality for the univariate case). For simplicity, we can assume that terms of second and higher order are small and be neglected. Eq. (4) then reduces to:

$$
\begin{equation*}
r_{i}=g_{i}\left[E\left(\delta_{j}\right)\right]+\sum_{j}\left[\delta_{j}-E\left(\delta_{j}\right)\right] \frac{\partial g_{i}}{\partial \delta_{j}}+\epsilon_{i} \tag{5}
\end{equation*}
$$

A stronger assumption involves the (local) linearity of $g_{i}$ in the $\delta_{j}$ 's. The sensitivity coefficients of the returns for the common factors are then constants $b_{i j}$ (for a specific range of $\delta_{j}$ ):

$$
\begin{equation*}
\frac{\partial r_{i}}{\partial \delta_{j}}=\frac{\partial g_{i}}{\partial \delta_{j}}=b_{i j} \quad, \quad \forall i \in N, \quad j \in K \tag{6}
\end{equation*}
$$

Under these assumptions, the RGP eq. (3) reduces to:

$$
\begin{equation*}
r_{i}=E\left(r_{i}\right)+\sum_{j}\left[\delta_{j}-E\left(\delta_{j}\right)\right] b_{i j}+\epsilon_{i}, \quad i \in N \tag{7}
\end{equation*}
$$

If the identities of the factors are known, this linear relation can be estimated by means of OLS regression. As easily can be seen from eqs. (5) and (7), the OLS estimation does not provide a Taylor series approximation. The argument is that the choice of the point of expansion is arbitrary (see eq. (4)).

White [1980] and Van Praag [1981], however, note that the OLS estimators $\left(b_{i j}\right)$ provide the best linear approximation of the dependent random variable (here: $r_{1}$ ) by the independent random variables (here: the $\delta_{j}^{\prime} s$ ): the $b_{i j}$ 's result from minimizing the mean square error of the approximation. For the least squares approximation, White assumes that there exists a functional model, although that may be unknown. Van Praag considers this regression approximation as 'model free', so that we even can interprete the sensitivities $b_{i j}$ without assuming any underlying functional model at all.

If there should exist a return generating model $g\left(\delta_{j}\right)_{j \in K}$ that is non-linear, it is important to note that problems can arise. As can be inferred from eqs. (5) and (7), the crucial point then is that the response coefficients $b_{i j}$ are only correct for small deviations of the factors. If the variability of the factors changes, this may imply that the sensitivities for these factors also change. In other words, the functional form eq. (7) then changes and does not apply any longer for different time periods (i.e. different levels of the factors). In this paper, however, we'll treat the factor sensitivities as constant technical coefficients; we'll model structural shifts in the economy as changes in the variances of the factors.

## 3. THE STATIONARITY OF $\beta$ WHEN RETURNS FOLLOW A LINEAR MULTI-FACTOR RGP

In this section, we investigate the implications of changing factor variances in a K-factor RGP for the $\beta$-coefficients of the securities. We start from the following assumptions:
(A1) For all securities, the RGP is linear in its parameters:

$$
\begin{equation*}
r_{i}=E\left(r_{i}\right)+\sum_{j} b_{i j} f_{j}+\epsilon_{i}, \quad \forall i \in N \tag{8}
\end{equation*}
$$

where $f_{j}$ denotes $\delta_{j}-E\left(\delta_{j}\right)$.
(A2) The common factors are mutually uncorrelated and uncorrelated with the idiosyncratic factors:

$$
\begin{equation*}
\operatorname{Cov}\left(f_{j}, f_{j},\right)=\operatorname{Cov}\left(f_{j}, \epsilon_{1}\right)=0, \quad \forall j \neq j^{\prime} \tag{9}
\end{equation*}
$$

(A3) The idiosyncratic factors are truly security specific, so

$$
\begin{equation*}
\operatorname{Cov}\left(\epsilon_{1}, \epsilon_{1},\right)=0, \quad \forall i \neq i^{\prime} \tag{10}
\end{equation*}
$$

As the sensitivities $b_{i j}$ link the (variability of the) security returns to the (variability of the) factors, we can say that these coefficients measure the factor risk of the securities. The variability of the error term $\epsilon_{1}$ induces the idiosyncratic risk.

The CAPM- $\beta$ links the (variability of the) security returns to the (variability of the) return on the market portfolio, so $\beta$ measures the market risk of the securities. As the market portfolio $m$ is a convex combination of the $N$ securities, its RGP can be expressed as:

$$
\begin{equation*}
r_{m}=\sum_{i} m_{i} r_{i}=E\left(r_{m}\right)+\sum_{j} b_{m j} f_{j}+\epsilon_{m} \tag{11}
\end{equation*}
$$

where $m_{i}$ is the proportion of the total market value of security $i$ relative to the market value of the market portfolio ( $\sum_{i} m_{i}=1$ ). Given the RGP with its assumptions, the covariances between the securities' returns and the return on the market portfolio can be expressed as:

$$
\begin{equation*}
\operatorname{Cov}\left(r_{i}, r_{m}\right)=\sum_{j} b_{i j} b_{m j} \operatorname{Var}\left(f_{j}\right)+m_{i} \operatorname{Var}\left(\epsilon_{i}\right), \forall i \in \mathbb{N} \tag{12}
\end{equation*}
$$

The variance of the return on the market portfolio equals:

$$
\begin{equation*}
\operatorname{Var}\left(r_{m}\right)=\sum_{j} b_{m j}{ }^{2} \operatorname{Var}\left(f_{j}\right)+\operatorname{Var}\left(\epsilon_{m}\right) \tag{13}
\end{equation*}
$$

Using eqs (12) and (13) in (1), we can rewrite $\beta_{i}$ in terms of factor variances:

$$
\begin{equation*}
\beta_{1}=\frac{\sum_{\mathrm{j}} \mathrm{~b}_{1 \mathrm{j}} \mathrm{~b}_{\mathrm{mj}} \operatorname{Var}\left(\mathrm{f}_{\mathrm{j}}\right)+\mathrm{m}_{\mathrm{i}} \operatorname{Var}\left(\epsilon_{\mathrm{i}}\right)}{\sum_{\mathrm{j}} \mathrm{~b}_{\mathrm{mj}}{ }^{2} \operatorname{Var}\left(\mathrm{f}_{\mathrm{j}}\right)+\operatorname{Var}\left(\epsilon_{\mathrm{m}}\right)}, \quad \forall i \in \mathrm{~N} \tag{14}
\end{equation*}
$$

The central issue is how $\beta$ behaves if the economic environment (in which the returns of the securities are generated) changes. For this purpose, we additionally assume that:
(A4) fundamental shifts in the structure of the economic environment arise because the variance of one or more common return generating factors changes.

Using eqs. (12) and (13), the changes in the relevant variancecovariance structure of the security returns can then be modelled as:

$$
\begin{equation*}
d \operatorname{Cov}\left(r_{i}, r_{m}\right)=\sum_{j} b_{i j} b_{m j} d \operatorname{Var}\left(f_{j}\right), \quad \forall i \in N \tag{15}
\end{equation*}
$$

and

$$
\begin{equation*}
d \operatorname{Var}\left(x_{\mathrm{m}}\right)=\sum_{\mathrm{j}} \mathrm{~b}_{\mathrm{mj}}^{2} d \operatorname{Var}\left(\mathrm{f}_{\mathrm{j}}\right) \tag{16}
\end{equation*}
$$

Given eqs. (15) and (16), the question arises how $\beta$ changes when the factor variances undergo small changes. If the changes in the numerator of $\beta$ in a way offset the changes in the denominator, $\beta$ will not change. But under what conditions will $\beta$ indeed be stationary?

To answer these questions, we have to analyze the change in $\beta$ as a result of a marginal change in a factor's variance. From eq. (14) it follows that:

$$
\begin{equation*}
\frac{\partial \beta_{i}}{\partial \operatorname{Var}\left(f_{j}\right)}=\frac{b_{i j} b_{m j} \operatorname{Var}\left(r_{m}\right)-\operatorname{Cov}\left(r_{i}, r_{m}\right) b_{m j}^{2}}{\left[\operatorname{Var}\left(r_{m}\right)\right]^{2}}, \quad \forall i \in N, j \in K \tag{17}
\end{equation*}
$$

### 3.1. Single-factor RGP

We first consider the case in which the RGP only consists of one factor, factor $k$, and where the variance of this factor exhibits a marginal change.

THEOREM 1: A security's $\beta$ will only be stationary under a marginal structural shift in the one-factor economy if this $\beta$ is equal to the security's factor sensitivity $b_{i k}$ scaled to the of the market portfolio's factor sensitivity $b_{m \mathrm{~m}}$

Proof: It will be clear that $\beta_{i}$ is invariant under a marginal change in $\operatorname{Var}\left(f_{k}\right)$ if:

$$
\begin{equation*}
\frac{\partial \beta_{\mathrm{i}}}{\partial \operatorname{Var}\left(\mathrm{f}_{\mathrm{k}}\right)}=0 \tag{18}
\end{equation*}
$$

Combining eqs. (17) and (18) and using the definition of $\beta_{i}$ yields the restriction

$$
\begin{equation*}
\beta_{i}=\frac{\mathrm{b}_{i \mathrm{k}}}{\mathrm{~b}_{\mathrm{mk}}} \tag{19}
\end{equation*}
$$

provided that the market portfolio has a non-zero sensitivity for the factor $k$. In a one-factor RGP, this is of course very likely.

The next two corollaries give insight in the conditions under which eq. (19) is satisfied.

COROLLARY 1: If the market portfolio is an equally weighted portfolio of a finite number of $N$ securities, the $\beta$-stationarity condition (19) requires the equality of the ratio of the security's factor risk to the security's idiosyncratic risk and the ratio of average factor risk to average idiosyncratic risk.

Proof: Substituting eq. (14) in (19), we get

$$
\begin{equation*}
\frac{b_{i k} b_{m k} \operatorname{Var}\left(f_{k}\right)+m_{1} \operatorname{Var}\left(\epsilon_{1}\right)}{b_{m k} 2 \operatorname{Var}\left(f_{k}\right)+\operatorname{Var}\left(\epsilon_{m}\right)}=\frac{b_{i k}}{b_{m k}} \tag{20}
\end{equation*}
$$

or

$$
\begin{equation*}
\frac{b_{i k}}{b_{m k}}=m_{i} \frac{\operatorname{Var}\left(\epsilon_{i}\right)}{\operatorname{Var}\left(\epsilon_{\mathrm{m}}\right)} \tag{21}
\end{equation*}
$$

When the market portfolio is equally weighted ( $m_{i}=1 / \mathrm{N}, \forall \mathrm{i}$ ), we have

$$
\begin{align*}
\operatorname{Var}\left(\epsilon_{\mathrm{m}}\right) & =\operatorname{Var}\left(\sum_{\mathrm{s} \in \mathrm{~m}} \mathrm{~m}_{\mathrm{s}} \epsilon_{\mathrm{s}}\right)=\operatorname{Var}\left(\mathrm{N}^{-1} \sum_{\mathrm{s}} \epsilon_{\mathrm{s}}\right)=\mathrm{N}^{-2} \sum_{\mathrm{s}} \operatorname{Var}\left(\epsilon_{\mathrm{s}}\right) \\
& =\frac{1}{\mathrm{~N}} \overline{\operatorname{Var}}(\epsilon) \tag{22}
\end{align*}
$$

where the bar over $\operatorname{Var}($.$) denotes its average value.$
Using this result, we can rewrite eq. (21) as

$$
\begin{equation*}
\frac{b_{i k}}{\operatorname{Var}\left(\epsilon_{i}\right)}=\frac{b_{m k}}{\overline{\operatorname{Var}}(\epsilon)} \tag{23}
\end{equation*}
$$

For a typical security, the idiosyncratic variance of its return equals the average of all securities' idiosyncratic variances. For its $\beta$ to be stationary, its sensitivity to factor $k$ must equal the market average sensitivity $b_{m k}$ to this factor. Consequently (according to eq. (19)), its $\beta$ must equal one. For a security with a larger (smaller) than typical idiosyncratic variance, $b_{1 k}$ must be larger (smaller) than $b_{m k}$, which implies that its $\beta$ must be larger (smaller) than one.

COROLLARY 2: If the market portfolio consists of a (countably) infinite number of securities, each with an infinitesimally small weight, then all $\beta^{\prime} s$ will be stationary in the one-factor economy.

Proof: If the number of securities in the market portfolio increases more and more, it follows from eq. (22) that the market portfolio
becomes well-diversified with respect to the return generating factor:
$\operatorname{Var}\left(\epsilon_{\mathrm{m}}\right) \rightarrow 0$ as $\mathrm{N} \rightarrow \infty$. Without variability there can exist no covariability, so $\operatorname{Cov}\left(\epsilon_{\mathrm{i}}, \epsilon_{\mathrm{m}}\right)=(1 / \mathrm{N}) \operatorname{Var}\left(\epsilon_{\mathrm{i}}\right) \rightarrow 0$ likewise. As a result of this naive diversification process, we can rewrite eq. (20) as

$$
\begin{equation*}
\frac{b_{i k} b_{m k} \operatorname{Var}\left(f_{k}\right)}{b_{m k}{ }^{2} \operatorname{Var}\left(f_{k}\right)}=\frac{b_{i k}}{b_{m k}} \tag{24}
\end{equation*}
$$

Hence, the stationarity condition (19) is satisfied by any value of $b_{i k}$ and $b_{m k}$.

Note that the $\beta$-stationarity conditions, as stated in the corollaries above, are independent of $\operatorname{Var}\left(f_{k}\right)$. Consequently, these conditions do not only apply for marginal changes in the factor's variance, but also for any arbitrary large changes.

### 3.2. Multi-factor RGP

Empirical evidence (as early as King [1966]) indicates that a single-factor RGP is an oversimplification of reality and that a multifactor RGP is more appropriate. In this general case, the economic environment is allowed to change in more than one dimension, i.e. the variances of K factors are allowed to change ( $\mathrm{K}>1$ ).

THEOREM 2: A security's $\beta$ will only be stationary under any marginal structural shift in the multi-factor economy iff (i) the security's factor sensitivities $\left(b_{i j}\right\}_{j \in K}$ are (pairwise) equal to the market portfolio's factor sensitivities $\left(b_{m \mathrm{j}}\right)_{\mathrm{j} \in \mathrm{K}}$, and (ii) its $\beta$ equals one.

Proof: The condition that a change in $\beta$ as a result of marginal changes in the factor variances is zero can be expressed as:

$$
\begin{equation*}
d \beta_{i}=\sum_{j} \frac{\partial \beta_{i}}{\partial \operatorname{Var}\left(f_{j}\right)} d \operatorname{Var}\left(f_{j}\right)=0 \tag{25}
\end{equation*}
$$

Using eq. (17), this is equivalent to

$$
\begin{equation*}
\frac{1}{\operatorname{Var}\left(\mathrm{r}_{\mathrm{m}}\right)} \sum_{\mathrm{j}}\left[\mathrm{~b}_{\mathrm{ij}} \mathrm{~b}_{\mathrm{mj}}-\beta_{\mathrm{i}} \mathrm{~b}_{\mathrm{mj}}^{2}\right] d \operatorname{Var}\left(\mathrm{f}_{\mathrm{j}}\right)=0 \tag{26}
\end{equation*}
$$

Solving for $\beta_{1}$, this yields

$$
\begin{equation*}
\beta_{i}=\frac{\sum_{j} b_{i j} b_{m j} d \operatorname{Var}\left(f_{j}\right)}{\sum_{j} b_{m j}^{2} d \operatorname{Var}\left(f_{j}\right)} \tag{27}
\end{equation*}
$$

It follows that $\beta_{i}$ is only stationary under independent marginal changes in the variances of the return generating factors if
(i) $\quad b_{i j}=b_{m j}, \forall j \in K$

According to eq. (27), we then also must have
(ii) $\quad \beta_{i}=1$. ■

The derivation of these conditions from eq. (27) does not depend on the magnitude of the factor variance changes. Consequently, these conditions apply for any arbitrary large changes in the factor variances.

Given eq. (14), the stationarity conditions (i) and (ii) imply

$$
\begin{equation*}
\mathrm{m}_{\mathrm{i}} \operatorname{Var}\left(\epsilon_{\mathrm{i}}\right)=\operatorname{Var}\left(\epsilon_{\mathrm{m}}\right) \tag{28}
\end{equation*}
$$

If the market portfolio is an equally weighted portfolio of a finite number of securities, eq. (22) applies. Consequently, only typical securities can satisfy these stationarity conditions (cf. corollary 1). If, on the other hand, the market portfolio consists of a (countably) infinite number of securities, each with an infinitesimally small weight, then idiosyncratic risk is no longer relevant. Condition (ii) becomes redundant and the pairwise equality of the security's and market
portfolio's factor sensitivities then forms the necessary and sufficient stationarity condition.

Condition (i) places a severe restriction on the variancecovariance structure of the security returns. If this condition is satisfied, the covariances between the returns of all "intertemporal stationary $\beta^{\prime \prime}$ securities would be identical. In a large, well-diversified capital market, these covariances would equal the variance of the return on the market portfolio. Also, the variances of the returns on these individual securities must then all be larger than the variance of the market portfolio's return (this follows from eq. (13): the market portfolio is well-diversified, but the security returns can contain an idiosyncratic component).

## 4. CONCLUSIONS

We hypothesize that fundamental shifts in the structure of the uncertain economic environment, in which the returns of the securities are generated, arise because the variance of one or more common underlying return generating factors changes. Using this return generating process, by means of which changes in the relevant covariance structure of the security returns are modelled, we analyzed the conditions under which a security's $\beta$-coefficient will be stationary.

In a single-factor context, all $\beta^{\prime} s$ will be stationary if the market portfolio is large, equally weighted and hence well-diversified with respect to the return generating factor. Should the market portfolio in contrast contain idiosyncratic risk, then for a security- $\beta$ to be stationary there should exist a restrictive relationship between the security's factor risk and idiosyncratic risk. More specific, for that security the ratio factor risk to average (market) factor risk must equal the ratio idiosyncratic risk to average idiosyncratic risk.

In a more realistic multi-factor context, security $\beta^{\prime} s$ will only be stationary if the large, equally weighted market portfolio is well-diversified with respect to the factors and the factor risks of the securities are pairwise equal to the factor risks of the market portfolio. This implies a severe restriction on the variance-covariance structure
of security returns, which is not likely to be satisfied in reality. In case the market portfolio should contain idiosyncratic risk, we have the additional stationarity condition that $\beta^{\prime} s$ must equal one. Hence, only securities with average factor risks and average idiosyncratic risk would exhibit stationary $\beta^{\prime} s$.

The relevance of results is two-fold. From an ex ante point of view, the ability to predict changes in estimated $\beta^{\prime}$ s could be enhanced by empirically identifying the factors that generate the returns and thus contribute to $\beta$-nonstationarity. The ability to identify factors contributing to intertemporal changes in $\beta$ also has potential value in an ex post context. Event studies rely on constant- $\beta$ market models to estimate residual returns. If $\beta$ changes over time, an error is made in estimating the residuals. It would be more correct to adjust the residuals for the changes in $\beta$. In the evaluation of investment performance, errors may result when relying on constant- $\beta$ market models; adjustments of $\beta$ according to changes in the economic environment would be adequate. Although empirical research has yielded some candidates for the factors (cf. Chen, Roll \& Ross [1986]), it would be worthwile to investigate the intertemporal behavior of their variability and link these results to the behavior of $\beta$.

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# APPLICATION OF NEUTS' METHOD <br> TO VACATION MODELS <br> WITH BULK ARRIVALS 

## Matendo KAYEMBE

Université de Mons-Hainaut
Place Warocqué 17, B-7000 Mons, Belgium

## ABSTRACT

This paper studies a single server infinite capacity queueing system with Poisson arrivals of customers groups of random size and a general service time distribution, the server of which applies a general exhaustive service vacation policy. A computational method is applied to obtain the steady state distributions of the queue of a post-departure or inactive phase termination epoch, at a post-departure epoch and at an arbitrary epoch. Relations between these distributions are given. As special cases, we consider two hybrid vacation models: the ( $T(S V), N$ )-policy and the ( $T(M V$ ); $N$ )-policy. In particular, when the service time distribution is of phase type, explicit results can be obtained for the N -policy. Whe show this by a simple example.

## 1. Introduction

This paper deals with the single server infinite capacity queueing system - considered in Loris-Teghem (1990a) --satisfying the following hypotheses:
$H_{1}$ : groups of customers arrive according to a Poisson process of parameter $\lambda$. The sizes of the groups are i.i.d. random variables, with probability distribution $\left\{d_{k}\right\}_{k \geqslant 1}$, generating function $D(z) \equiv \sum_{k \geqslant 1} d_{k} z^{k}(|z| \leqslant 1)$, finite expectation $\mathbb{E}[D]$ and second order moment $\mathbb{E}\left[D^{2}\right]$.
$H_{2}$ : The server alternates between active and inactive states. In the active state, the server provides service to customers, so that in what we call an "active phase", the system is never empty. The epochs $0 \leqslant t_{0}<t_{1}<\cdots<t_{m}<\cdots$ at which the server "enters" the inactive state - and thus, becomes unavailable to the customers - are the times at which the system gets empty (exhaustive service). We denote by $\tau_{m}$, $m \geqslant 1$, the epoch at which the inactive phase beginning at $t_{m-1}$ terminates. Thus, if $X_{m}, m \geqslant 1$, is the number of groups of customers arriving in the time interval $\left.] t_{m-1}, \tau_{m}\right]$, we have $X_{m} \geqslant 1$ a.s. We put $V_{m}=\tau_{m}-t_{m-1}, m \geqslant 1$, and we assume that the random variables $V_{m}, m \geqslant 1$, are i.i.d., with finite expectation and with distribution function $V(\cdot)$. We denote by $(X, V)$ a random vector distributed as the $\left(X_{m}, V_{m}\right)(m \geqslant 1)$.
$H_{3}$ : Customers are served individually in an order independent of their service times, which are i.i.d. random variables independent of the arrival process and the sequence $\left\{V_{m}\right\}_{m \geqslant 1}$, with distribution function $S(\cdot)$, Laplace-Stieltjes transform (L.S.T.) $\tilde{S}(\cdot)$, and with finite positive expectation $\mathbb{E}[S]$ and second order moment $\mathbb{E}\left[S^{2}\right]$.
$H_{\mathbf{4}}: \rho=\lambda \mathbb{E}[D] \mathbb{E}[S]<1$.
For this model, Loris-Teghem (1990a, 1990b) gives, in terms of generating functions, a probabilistic proof of the stochastic decomposition property for the stationary queue length distribution, at a post-departure epoch and at an arbitrary epoch.

Our purpose here is to obtain computational expressions for the stationary queue length distributions.
Computable formulas for a single server queue with server vacations are given in Lu cantoni, Meier-Hellstern and Neuts (1990). In the model considered in that paper, customers arrive (individually) according to a general arrival process, the Markovian Arrival Process - of which the Poisson Arrival Process is a very special case - and the vacation policy applied by the server is the "multiple vacation policy", sometimes called the "T(MV)-policy". The computable expressions obtained by the authors concern waiting
time distributions as well as queue length distributions, namely the stationary distribution of the queue length at a post-departure epoch, at an arrival epoch and at an arbitrary epoch.

In the present paper, we assume that groups of customers arrive according to a Poisson process. However, we consider a class of exhaustive service vacation policies which contains the $T(M V)$-policy as a special case. We are interested in the queue length process, and we first consider the queue length at service completion or inactive phase termination epochs (section 2), for which we compute the steady-state distribution by Neuts'method. We then relate the steady-state distribution of the queue length at service completion epochs (section 3) and at an arbitrary epoch (section 4) to the former distribution. Section 5 is devoted to two particular policies: the ( $T(S V) ; N$ )-policy and the $(T(M V) ; N)$-policy. For $N=1$, the latter reduces to the $T(M V)$-policy considered in Lucantoni, Meier-Hellstern and Neuts. Our results for this policy agree with those obtained by particularizing the queue length results in Lucantoni, Meier-Hellstern, and Neuts to the Poisson Arrival Process

We mention that a general batch arrival process, the Batch Markovian Arrival Process, has been considered recently by Lucantoni (1991). We plan to further investigate vacation models, by considering such an arrival process.

## Notations

Let $Y_{m}, m \geqslant 1$, be the number of customers arriving in the time interval $\left.] t_{m-1}, \tau_{m}\right]$; $\left\{Y_{m}\right\}_{m \geqslant 1}$ is a sequence of i.i.d. random variables with $Y_{m} \geqslant 1$ a.s. We denote by $Y$ a random variable distributed as the $Y_{m}, m \geqslant 1$, with finite expectation $\mathbb{E}[Y]$ and second order moment $\mathbb{E}\left[Y^{2}\right]$.
For $t \geqslant 0, n \geqslant 0, \quad \operatorname{Re}(s) \geqslant 0$ and $|z| \leqslant 1$, define:

$$
\begin{aligned}
P(n, t) & =\operatorname{Pr}[n \text { customers arrive in the time interval }] 0, t]] \\
& =\sum_{k=0}^{n} e^{-\lambda t} \frac{(\lambda t)^{k}}{k!} d_{n}(k)
\end{aligned}
$$

where $\left\{d_{n}(k)\right\}_{n \geqslant 0}$ is the $k$-fold convolution of the probability distribution $\left\{d_{n}\right\}_{n \geqslant 1}$.

$$
\begin{align*}
A_{n}(t) & =\int_{0}^{t} P(n, x) d S(x)  \tag{1}\\
\dot{A}_{n}(s) & =\int_{0}^{\infty} e^{-s x} d A_{n}(x)=\int_{0}^{\infty} e^{-s x} P(n, x) d S(x) \\
\stackrel{\circ}{A}(t) & =\sum_{n \geqslant 0} A_{n}(t)=S(t)
\end{align*}
$$

$$
\begin{align*}
\dot{A}(z, s) & =\sum_{n \geqslant 0} \dot{A}_{n}(s) z^{n}=\dot{S}[s+\lambda-\lambda D(z)] \\
A_{n} & =A_{n}(+\infty)=\dot{A}_{n}(0)=\sum_{j=0}^{n} \psi_{j} d_{n}(j) \tag{2}
\end{align*}
$$

where $\psi_{j}=\int_{0}^{\infty} e^{-\lambda t} \frac{\left.(\lambda t)^{j}\right)}{j!} d S(t), j \geqslant 0$

$$
\begin{align*}
A(z) & =\sum_{n \geqslant 0} A_{n} \tilde{z}^{n}=\tilde{A}(z, 0)=\check{S}[\lambda-\lambda D(z)]  \tag{3}\\
A_{n}^{\star} & =\int_{0}^{\infty} P(n, t)[1-S(t)] d t=\sum_{j=0}^{n} \psi_{j}^{*} d_{n}(j) \tag{4}
\end{align*}
$$

where $\psi_{j}^{*}=\int_{0}^{\infty} e^{-\lambda t\left(\frac{(\lambda t)}{j!}\right.}[1-S(t)] d t, j \geqslant 0$.

$$
\begin{equation*}
C_{n}(t)=\operatorname{Pr}[V \leqslant t, Y=n] \tag{5}
\end{equation*}
$$

$$
\begin{align*}
& \left(C_{0}(t) \equiv 0\right) \\
\dot{C}_{n}(s)= & \int_{0}^{\infty} e^{-s t} d C_{n}(t) \\
\dot{C}(t)= & \sum_{n \geqslant 1} C_{n}(t)=V(t) \\
\dot{C}(z, s)= & \sum_{n \geqslant 1} \dot{C}_{n}(s) z^{n} \\
C_{n}= & C_{n}(+\infty)=\dot{C}_{n}(0)=\operatorname{Pr}[Y=n]  \tag{6}\\
C(z)= & \sum_{n \geqslant 1} C_{n} z^{n}=\dot{C}(z, 0)
\end{align*}
$$

## Remark 1

(a) For $n \geqslant 1 . C_{n}(t)$ depends on the specific type of vacation policy considered.
(b) Explicit expressions can be obtained for $\psi_{j}$ and $\psi_{j}^{*}, j \geqslant 0$, when the service time distribution $S(\cdot)$ is of phase type with the (irreducible) representation $(\underline{\beta}, H)$ of order $n$. Recall that this means that $S(\cdot)$ is the distribution of the time until absorption in a Markov process with a finite number $n$ of transient states $1,2, \ldots, n$ and a single absorbing state $n+1$. The infinitesimal generator of such a process is of the form

$$
\left(\begin{array}{cc}
H & H^{0} \\
0 & 0
\end{array}\right)
$$

where $H$ is a non singular $n \times n$ matrix which describes transitions between transient states, and $\underline{H}^{0}$ is a column vector of dimension $n$, which describes transitions from transient states to the absorbing state, with $\underline{H}^{0}=-H \underline{e}$. Throughout this paper, $\underline{e}$ is a column vector of appropriate dimension with all its components equal to $1 .\left(\underline{\beta}, \beta_{n+1}\right)$ is the initial probability vector and $S(x)$ is given by $S(x)=1-\underline{\beta} \exp (H x) \underline{e}$, for $x \geqslant 0$. For simplicity, we consider only the case where $\beta_{n+1}=0$.
Let

$$
L \otimes M=\left(\begin{array}{cccc}
L_{11} M & L_{12} M & \cdots & L_{1 k_{2}} M \\
\vdots & & & \\
L_{k_{1}} M & L_{k_{1}} M & \cdots & L_{k_{1} k_{2}} M
\end{array}\right)
$$

be the Kronecker product of two matrices $L$ and $M$ of dimensions $k_{1} \times k_{2}$ and $k_{1}^{\prime} \times k_{2}^{\prime}$ respectively ( $L \otimes M$ is a matrix of dimension $k_{1} k_{1}^{\prime} \times k_{2} k_{2}^{\prime}$ ). Let $(1,(-\lambda)$ ) denote the simplest representation of the group arrival Poisson process and let $I$ denote the identity matrix (of appropriate dimension). Then (see Neuts [1989, th. 5.1.5, pp. 244-246])

$$
\begin{aligned}
& \psi_{0}^{\prime}=-(I \otimes \underline{\beta})((-\lambda) \otimes I+I \otimes H)^{-1}\left(I \otimes \underline{H}^{0}\right) \\
& \psi_{k}=\varphi \gamma^{k-1} \Omega, \quad k \geqslant 1,
\end{aligned}
$$

and

$$
\begin{align*}
& \psi_{0}^{*}=-(I \ominus \underline{\beta})((-\lambda) \otimes I+I \otimes H)^{-1}(I \otimes \underline{\varepsilon})  \tag{8}\\
& \psi_{k}^{*}=\varphi \gamma^{k-1} \Omega_{0}, \quad k \geqslant 1
\end{align*}
$$

where the matrices $\varphi, \gamma, \Omega$ and $\Omega_{0}$ (of dimensions $1 \times n, n \times n$ and $n \times 1$ respectively) are given by

$$
\begin{align*}
\varphi & =-(I \otimes \underline{\beta})((-\lambda) \otimes I+I \otimes H)^{-1}(\lambda \otimes I) \\
\gamma & =-(1 \otimes I)((-\lambda) \otimes I+I \otimes H)^{-1}(\lambda \otimes I)  \tag{9}\\
\Omega & =-(1 \otimes I)((-\lambda) \otimes I+I \otimes H)^{-1}\left(I \otimes \underline{H}^{0}\right) \\
\Omega_{0} & =-(1 \otimes I)((-\lambda) \otimes I+I \otimes H)^{-1}(I \otimes \underline{e})
\end{align*}
$$

From (7) and (8), it is clear that no numerical integrations are needed to compute $\psi_{k}$ and $\psi_{k}^{\star}, k \geqslant 0$. We refer to Latouche (1989) and Neuts (1989) for more information about
phase type distributions

## 2. The queue length at service completion or inactive phase termination

 epochsLet $N(t)$ be the number of customers in the system at time $t$. We are concerned with the discrete parameter process obtained by restricting the process $\{N(t), t \geqslant 0\}$ at epochs of service completion or inactive phase termination. Denote these instants by $\theta_{n}, n \geqslant 0$, and let $\zeta_{n}, n \geqslant 0$, be the number of customers arrived in the time interval $] \theta_{n}, \theta_{n+1}$ ] Define

$$
T_{n}=\theta_{n+1}-\theta_{n} \quad \text { and } \quad N_{n}=N\left(\theta_{n}+\right), n \geqslant 0
$$

Then, for $n \geqslant 0$,

$$
N_{n+1}=\left(N_{n}-1\right)^{+}+\zeta_{n}
$$

and

$$
\begin{aligned}
& \operatorname{Pr}\left[N_{n+1}=j, T_{n} \leqslant x \mid N_{0}, T_{0}, N_{1}, T_{1}, \ldots, T_{n-1}, N_{n}=i\right] \\
& \left.=\operatorname{Pr}\left[N_{n+1}=j, T_{n} \leqslant x \mid N_{n}=i\right] \quad\left(\equiv Q_{i j}(x)\right), x \geqslant 0\right)
\end{aligned}
$$

It follows that the sequence $\left\{\left(N_{n}, T_{n}\right)\right\}_{n \geqslant 0}$ is a M.R.P. on the state space $\{i \geqslant 0\} \times[0,+\infty[$ Its transition probability matrix $Q(\cdot)$ is given by

$$
Q(x)=\left(\begin{array}{cccc}
C_{0}(x) & C_{1}(x) & C_{2}(x) & \cdots  \tag{10}\\
A_{0}(x) & A_{1}(x) & A_{2}(x) & \cdots \\
0 & A_{0}(x) & A_{1}(x) & \cdots \\
0 & 0 & A_{0}(x) & \cdots \\
\vdots & \vdots & \vdots &
\end{array}\right), x \geqslant 0
$$

where $A_{\nu}(x)$ and $C_{\nu}(x), \nu \geqslant 0$ are the probability mass-functions defined in (1) and (5). Note that

$$
\begin{aligned}
\rho & =\sum_{k \geqslant 1} k A_{k}<1 \quad\left(H_{4}\right), \\
\stackrel{\circ}{\alpha} & \equiv \int_{0}^{\infty} x d \stackrel{\circ}{A}(x)=\int_{0}^{\infty} x d S(x)=\mathbb{E}[S]<+\infty, \\
\stackrel{\circ}{\beta} & \equiv \int_{0}^{\infty} x d \stackrel{\circ}{C}(x)=\int_{0}^{\infty} x d V(x)=\mathbb{E}[V]<+\infty, \\
\sum_{k \geqslant 1} k C_{k} & =C^{\prime}(1)=\mathbb{E}[Y]<+\infty .
\end{aligned}
$$

Thus (see Neuts [1989, th. 1.3.1, pp. 12-13]), the M.R.P. $\left\{\left(N_{n}, T_{n}\right)\right\}_{n \geqslant 0}$ and the Markov chain $\left\{N_{n}\right\}_{n \geqslant 0}$ are positive recurrent. From (10), the transition probability matrix of the Markov chairı $\left\{N_{n}\right\}_{n \geqslant 0}$ is given by

$$
Q=Q(+\infty)=\left(\begin{array}{cccc}
C_{0} & C_{1} & C_{2} & \cdots  \tag{11}\\
A_{0} & A_{1} & A_{2} & \cdots \\
0 & A_{0} & A_{1} & \cdots \\
0 & 0 & A_{0} & \cdots \\
\vdots & \vdots & \vdots &
\end{array}\right)
$$

Clearly, this Markov chain is irreducible and aperiodic.
Let $\underline{x}=\left(x_{i}\right)_{i \geqslant 0}$ be the stationary probability vector, which satisfies

$$
\begin{aligned}
\underline{x} Q & =\underline{x} \\
\underline{x} \underline{e} & =1
\end{aligned}
$$

From (11), we have

$$
\begin{equation*}
x_{i}=x_{0} C_{i}+\sum_{\nu=1}^{i+1} x_{\nu} A_{i+1-\nu}, \text { for } i \geqslant 0 \tag{12}
\end{equation*}
$$

2.1. Computation of $x_{0}$

From Neuts (1989, p. 15) we obtain

$$
\begin{equation*}
x_{0}=\left[1+\frac{1}{1-\rho} \sum_{\nu \geqslant 1} \nu C_{\nu}\right]^{-1}=\frac{1-\rho}{1-\rho+\mathbb{E}[Y]} \tag{13}
\end{equation*}
$$

2.2. Computation of $x_{i}, i \geqslant 1$

From Neuts (1989, p. 17) we have the following recurrence formula:

$$
\begin{equation*}
x_{i}=A_{0}^{-1}\left[x_{0} \hat{C}_{i-1}+\sum_{\nu=1}^{i-1} x_{\nu} \hat{A}_{i-\nu}\right], i \geqslant 1 \tag{14}
\end{equation*}
$$

where $\hat{A}_{\nu}=1-\sum_{r=0}^{\nu} A_{r}$, and $\hat{C}_{\nu}=1-\sum_{r=0}^{\nu} C_{r}$, for $\nu \geqslant 0$.
An analogous recurrence formula for ${ }^{r=0}$ the computation of $x_{i}, i \geqslant 1$, can be found in Ramaswami (1988).
2.3. Mean queue length at service completion or inactive phase termination epochs.

Put $X_{v}(z)=\sum_{i \geqslant 0} x_{i} z^{i}(|z| \leqslant 1)$. Equations (12) readily imply

$$
X_{v}(z)[z-A(z)]=x_{0}[z C(z)-A(z)]
$$

so that (see Neuts [1989, th. 1.4.1, pp. 17-18]) $\sum_{i \geqslant 1} i x_{i} \equiv X_{v}^{\prime}(1)$ is given by

$$
\begin{equation*}
X_{v}^{\prime}(1)=\frac{1}{2(1-\rho)}\left\{A^{\prime \prime}(1)+x_{0}\left[2 C^{\prime}(1)+C^{\prime \prime}(1)-A^{\prime \prime}(1)\right]\right\} \tag{15}
\end{equation*}
$$

where $x_{0}$ is obtained in (13),

$$
A^{\prime \prime}(1)=\sum_{\nu \geqslant 2} \nu(\nu-1) A_{\nu}=\lambda^{2}(\mathbb{E}[D])^{2} \mathbb{E}\left[S^{2}\right]+\lambda \mathbb{E}[S]\left(\mathbb{E}\left[D^{2}\right]-\mathbb{E}[D]\right)
$$

and

$$
C^{\prime \prime}(1)=\sum_{\nu \geqslant 2} \nu(\nu-1) C_{\nu}=\mathbb{E}\left[Y^{2}\right]-\mathbb{E}[Y] .
$$

## 3. The queue length at post-departure epochs

By restricting the process $\{N(t), t \geqslant 0\}$ at service completion epochs, we obtain another M.R.P. on the state space $\{i \geqslant 0\} \times[0,+\infty[$, with transition probability matrix $Q_{1}(\cdot)$ given by

$$
Q_{1}(x)=\left(\begin{array}{cccc}
B_{0}(x) & B_{1}(x) & B_{2}(x) & \cdots  \tag{16}\\
A_{0}(x) & A_{1}(x) & A_{2}(x) & \cdots \\
0 & A_{0}(x) & A_{1}(x) & \cdots \\
0 & 0 & A_{0}(x) & \cdots \\
\vdots & \vdots & \vdots &
\end{array}\right), x \geqslant 0
$$

where $B_{\mathbf{i}}(x)=\sum_{\nu=1}^{i+1} C_{\nu}(\cdot) \star A_{i+1-\nu}(x)$, for $i \geqslant 0$.
Thus $Q_{1}(x)$ differs from $Q(x)$-given by (10)-only in the elements on the first line: $B_{i}(x)$ instead of $C_{i}(x)$, and the stationary probability vector $\underline{\pi}=\left(\pi_{i}\right)_{i \geqslant 0}$ of the corresponding Markov chain could be computed in a way similar to that described in section 2 for the computation of the stationary probability vector $\underline{x}$.

Note, however, that a simple probabilistic argument leads to the following relation

$$
\begin{equation*}
x_{i}=\frac{x_{0}}{\pi_{0}} \pi_{i}+x_{0} C_{i}, \quad i \geqslant 1 \tag{17}
\end{equation*}
$$

Indeed, as the system is not empty at an inactive phase termination, we have:

$$
\begin{aligned}
x_{0}= & \operatorname{Pr}[\text { the system is empty after a transition }] \\
= & \operatorname{Pr}[\text { the system is empty after a transition and this } \\
& \text { transition is a service completion }]
\end{aligned}
$$

$=\operatorname{Pr}[$ a transition is a service completion $] \times$ $\times \operatorname{Pr}$ [the system is empty after a transition, given that this transition is a service completion]
$=\operatorname{Pr}[$ a transition is a service completion $] \times \pi_{0}$
Thus, $\frac{x_{0}}{\pi_{0}}$ is the probability that a transition is a service completion.
On the other hand,

$$
\begin{aligned}
& \operatorname{Pr}[\text { a transition is an inactive phase termination }] \\
& =\operatorname{Pr}[\text { the system is empty after the previous transition }] \\
& =x_{0}
\end{aligned}
$$

The result follows by application of the law of total probability, as $C_{\mathrm{i}}$ is the probability that $i$ customers arrive in an inactive phase.
From (17) we have

$$
\sum_{i \geqslant 1} x_{i}=\frac{x_{0}}{\pi_{0}} \sum_{i \geqslant 1} \pi_{i}+x_{0} \sum_{i \geqslant 1} C_{i}
$$

or

$$
1-x_{0}=\frac{x_{0}}{\pi_{0}}\left(1-\pi_{0}\right)+x_{0}
$$

so that

$$
\begin{align*}
\pi_{0} & =\frac{x_{0}}{1-x_{0}}=\frac{1-\rho}{\mathbb{E}[Y]}  \tag{18}\\
\pi_{i} & =\frac{1}{1-x_{0}}\left(x_{i}-x_{0} C_{i}\right), i \geqslant 1
\end{align*}
$$

and, using (15),

$$
\begin{equation*}
\sum_{i \geqslant 1} i \pi_{i}=\rho+\frac{A^{\prime \prime}(1)}{2(1-\rho)}+\frac{C^{\prime \prime}(1)}{2 \mathbb{E}[Y]} \tag{19}
\end{equation*}
$$

## 4. The stationary queue length at an arbitrary epoch

In this section, we are interested in computing the continuous time stationary queue length distribution, denoted by $\left\{p_{i}\right\}_{i \geqslant 0}$, by using one of the M.R.P. considered in the previous sections. As we noticed, the matrices $Q(\cdot)$ and $Q_{1}(\cdot)$ associated with these processes differ only by the elements on the first line. As the expression for these elements is simpler in $Q(\cdot)$ than in $Q_{1}(\cdot)$, we will relate the distribution $\left\{p_{i}\right\}_{i \geqslant 0}$ to the stationary distribution $\left\{x_{i}\right\}_{i \geqslant 0}$ of the queue length at a service completion or inactive phase termination epoch. (In Loris-Teghem (1990a, 1990b), a relation between the distributions $\left\{p_{i}\right\}_{i \geqslant 0}$ and $\left\{\pi_{i}\right\}_{i \geqslant 0}$ is obtained).

We first introduce the fundamental mean $E$ of the M.R.P. $\left\{\left(N_{n}, T_{n}\right)\right\}_{n \geqslant 0}$ (i.e. the inner product of the stationary probability vector $\underline{x}$ and the vector of row sum means $\int_{0}^{\infty} x d Q(x) \underline{e}$ of the matrix $\left.Q(\cdot)\right)$. In the stationary version of the queue, $E$ may be interpreted as the average time between two consecutive transitions (i.e. service completion or inactive phase termination). From Neuts (1989, p. 22) we have

$$
\begin{align*}
E & =x_{0} \dot{\circ}+\left(1-x_{0}\right) \stackrel{\circ}{\alpha}=\frac{\stackrel{\circ}{\beta}(1-\rho)+\stackrel{\circ}{\alpha} \sum_{\nu \geqslant 1} \nu C_{\nu}}{1-\rho+\sum_{\nu \geqslant 1} \nu C_{\nu}} \\
& =\frac{\lambda^{-1} \mathbb{E}[X]-\mathbb{E}[S](\mathbb{E}[D] \mathbb{E}[X]-\mathbb{E}[Y])}{1-\rho+\mathbb{E}[Y]} \tag{20}
\end{align*}
$$

We assume that time $t=0$ corresponds to a transition epoch in the M.R.P. $\left\{\left(N_{n}, T_{n}\right)\right\}_{n \geqslant 0}$ and that $N_{0}=i_{0} \geqslant 0$.
For $t \geqslant 0$, let $M_{i_{0} i}(t)$ denote the conditional expected number of visits to state $i$ in the time interval $[0, t]$, given that $N_{0}=i_{0}$, and let $M(t)=\left\{M_{i_{0} i}(t), i_{0} \geqslant 0, i \geqslant 0\right\}$ be the Markov renewal matrix of the M.R.P. $\left\{\left(N_{n}, T_{n}\right)\right\}_{n \geqslant 0}$.
The matrix $M(\cdot)$ is given by the convolution power series

$$
M(t)=\sum_{n \geqslant 0} Q^{(n)}(t), \text { for } t \geqslant 0
$$

Let $d M_{i_{0} i}(t)$ denote the conditional probability that in the interval $] t, t+d t[$, the M.R.P. $\left\{\left(N_{n}, T_{n}\right)\right\}_{n \geqslant 0}$ enters state $i$, given that $N_{0}=i_{0}$.
We now consider the continuous-parameter process $\{N(t), t \geqslant 0\}$ and define:

$$
P_{i_{0} i}(t)=\operatorname{Pr}\left[N(t)=i \mid N_{0}=i_{0}\right], \text { for } i \geqslant 0
$$

and

$$
K_{i_{0} i}(t)=\operatorname{Pr}\left[N(t)=i, \theta_{1}>t \mid N_{0}=i_{0}\right], \text { for } i \geqslant i_{0}
$$

We have (see Çinlar $(1969,1975)$ )

$$
\begin{aligned}
P_{i_{0} i}(t) & =\sum_{j=0}^{i} \int_{0}^{t} d M_{i_{0} j}(u) K_{j i}(t-u) \\
& = \begin{cases}\int_{0}^{t} d M_{i_{0} 0}(u) e^{-\lambda(t-u)} ; & \text { for } i=0 \\
\int_{0}^{t} d M_{i_{0} 0}(u) K_{0 i}(t-u)+\sum_{j=1}^{i} \int_{0}^{t} d M_{i_{0} j}(u)[1-S(t-u)] P(i-j, t-u), \\
\text { for } i \geqslant 1 .\end{cases}
\end{aligned}
$$

The limits $p_{i}=\lim _{t \rightarrow \infty} P_{i_{0} i}(t), i \geqslant 0$, exist (and are independent of $i_{0}$ ) by virtue of the key renewal theorem and are given by

$$
\begin{align*}
p_{0} & =\frac{x_{0}}{E} \int_{0}^{\infty} K_{00}(t) d t=\frac{x_{0}}{\lambda E} \\
& =\frac{\lambda^{-1}(1-\rho)}{\lambda^{-1} \mathbb{E}[X]-\mathbb{E}[S](\mathbb{E}[D] \mathbb{E}[X]-\mathbb{E}[Y])},  \tag{21}\\
p_{i} & =\frac{1}{E} x_{0} \int_{0}^{\infty} K_{0 i}(t) d t+\frac{1}{E} \sum_{j=1}^{i} x_{j} \int_{0}^{\infty} P(i-j, t)[1-S(t)] d t, i \geqslant 1, \tag{22}
\end{align*}
$$

provided that each of the functions $K_{0 i}(\cdot)$ - which depends on the specific type of vacation policy considered - be directly Riemann integrable.

## Mean queue length at an arbitrary epoch

Relations (21) and (22) give, for the generating function $P_{v}(z) \equiv \sum_{i \geqslant 0} p_{i} z^{i}(|z| \leqslant 1)$ :

$$
\begin{equation*}
P_{v}(z)=p_{0}+\frac{1}{E}\left\{x_{0} \int_{0}^{\infty} K_{0}(t, z) d t+\left[X_{v}(z)-x_{0}\right][1-A(z)][\lambda-\lambda D(z)]^{-1}\right\} \tag{23}
\end{equation*}
$$

where $K_{0}(t, z)=\sum_{i \geqslant 1} K_{0 i}(t) z^{i}$.
Using this expression for $P_{v}(z)$, one could obtain the mean queue length $\sum_{i \geqslant 1} i p_{i}=P_{v}^{\prime}(1)$.
Note, however, that this can be obtained by using the following decomposition property (Loris-Teghem (1990a, 1990b)):

$$
P_{v}(z)=P_{n v}(z) \Delta_{v}(z) \quad, \quad|z| \leqslant 1
$$

where $P_{n v}(z)$ is the generating function of the queve length at an arbitrary epoch in the bulk arrival model without vacations and $\Delta_{v}(z)$ is the following generating function:

$$
\Delta_{v}(z)=\frac{1-C(z)}{\mathbb{E}[X](1-D(z))}
$$

Thus

$$
\begin{align*}
P_{v}^{\prime}(1) & =P_{n v}^{\prime}(1)+\Delta_{v}^{\prime}(1)  \tag{24}\\
\text { with } \quad P_{n v}^{\prime}(1) & =\rho+\frac{\lambda\left[\lambda(\mathbb{E}[D])^{2} \mathbb{E}\left[S^{2}\right]+\mathbb{E}[S]\left(\mathbb{E}\left[D^{2}\right]-\mathbb{E}[D]\right)\right]}{2(1-\rho)}
\end{align*}
$$

$\Delta_{v}^{\prime}(1)$ depends on the specific type of vacation policy considered.
5. Particular cases : the ( $T(S V) ; N)$-model and the $(T(M V) ; N)$-model

In this section, we are interested in the functions $K_{0 ;}(t), i \geqslant 1$, for two hybrid vacation models: the ( $T(S V) ; N)$-policy and the ( $T(M V) ; N$ )-policy, introduced in Loris-Teghem (1985). In these models, upon becoming idle, the server leaves the system for a vacation of random length. In the first case, when he comes back to the system, the server becomes active immediately if he finds at least $N$ customers present. Otherwise, he remains inactive, inspecting the queue, until $N$ customers are present. In the second case, the vacations are repeated until the server finds at least $N$ customers in the system upon return from a vacation. We refer the reader to Loris-Teghem (1990a) for more details about these policies. Note that in Loris-Teghem (1990b), the distribution $\left\{C_{k}\right\}_{k \geqslant 1}$ of $Y$ is given for both the ( $T(S V) ; N$ )-model and the ( $T(M V) ; N$ )-model.

Let $U(\cdot)$ be the common distribution function of the vacation lengths, with finite expectation $\mathbb{E}[U]$.
5.1. The ( $T(S V) ; N)$-model

We have :

$$
\begin{align*}
\kappa_{0 i}(t) & =P(i, t)[1-U(t)], & & \text { for } i \geqslant N  \tag{25}\\
& =P(i, t) & & \text { for } 1 \leqslant i \leqslant N-1
\end{align*}
$$

## Remark 2

For the $N$-policy (i.e. $U(\cdot) \equiv 1$ ), (25) reduces to

$$
\begin{aligned}
K_{0 i}(t) & =0 \quad, \text { for } i \geqslant N \\
& =P(i, t), \text { for } 1 \leqslant i \leqslant N-1
\end{aligned}
$$

Substituting into (22) yields

$$
\begin{align*}
p_{i} & =p_{0} \sum_{k=1}^{i} d_{i}(k)+\eta_{i} & , \text { for } 1 \leqslant i \leqslant N-1  \tag{26}\\
& =\eta_{i} & , \text { for } i \geqslant N
\end{align*}
$$

where $\eta_{i}(i \geqslant 1)$ is given by

$$
\eta_{i}=\frac{1}{E} \sum_{j=1}^{i} x_{j} A_{i-j}^{\star} \quad(i \geqslant 1)
$$

When the service time distribution is of phase type, explicit expressions can be obtained for $p_{i}(i \geqslant 1)$. This follows from (4) and remark 1 (b). In the following example, we write the expressions for $p_{0}, p_{1}, p_{2}$ and $p_{3}$.

Example
Consider the case where $n=N=2, H=\left(\begin{array}{cc}-\mu & \mu \\ 0 & -\mu\end{array}\right)$ and $\underline{\beta}=(1,0)$. We have

$$
\begin{aligned}
p_{0}= & \frac{1-\rho}{\mathbb{E}[X]}, \\
p_{1}= & p_{0}\left\{d_{1}+\frac{\lambda(\lambda+2 \mu)}{\mu^{2}}\right\}, \\
p_{2}= & \lambda p_{0}\left\{\frac{(\lambda+\mu)^{2}(\lambda+2 \mu)}{\mu^{4}}-\frac{\lambda}{\mu^{2}} d_{1}\right\}, \\
\text { and } \quad p_{3}= & \lambda p_{0}\left\{\frac{(\lambda+\mu)^{4}(\lambda+2 \mu)}{\mu^{6}}-\frac{(\lambda+\mu)\left(3 \lambda^{2}+5 \lambda \mu\right)}{\mu^{4}} d_{1}\right. \\
& \left.-\frac{(\lambda+2 \mu)}{\mu^{2}}\left(d_{1}^{2}+d_{2}\right)-\frac{\lambda}{\mu^{2}} d_{2}\right\},
\end{aligned}
$$

where $\mathbb{E}[X]=1+d_{1}$.

## Remark 3

The $N$-policy model with a phase type service time distribution was considered in Altiok (1987). In order to compute the $p_{i}, i \geqslant 1$, Altiok uses the following relations:

$$
\begin{align*}
p_{i} & =\sum_{k=0}^{n} P_{i, k}^{\star} \quad \text {, for } \quad 1 \leqslant i \leqslant N-1 \\
& =\sum_{k=1}^{n} P_{i, k}^{\star} \quad, \text { for } \quad i \geqslant N \tag{27}
\end{align*}
$$

where $P_{i, 0}^{\star}(0 \leqslant i \leqslant N-1)$ and $P_{i, k}^{\star}(i \geqslant 1, k=1, \ldots, n)$ are the steady-state probabilities that $i$ customers are in the system and the server is idle, and to have $i$ customers in the system and the server is in phase $k$ respectively, for the computation of which Altiok develops algorithms.

### 5.2. The $(T(M V) ; N)$-model

We have :

$$
\begin{align*}
\kappa_{0 i}(t) & =\sum_{r \geqslant 0} \sum_{k=0}^{N-1} \int_{0}^{t} P(k, y) d U^{(r)}(y) P(i-k, t-y)[1-U(t-y)], & , i \geqslant N  \tag{28}\\
& =P(i, t) & , 1 \leqslant i \leqslant N-1
\end{align*}
$$

where $U^{(r)}(t)$ denotes the $r$-th power of convolution of $U(t)$.

## Remark 4

Let $b_{0}=\int_{0}^{\infty} e^{-\lambda t} d U(t)$ denote the probability that no arrival occur during a single vacation.
Putting $N=1$ in

- (22) and (25) we get

$$
\begin{equation*}
p_{i}=\lambda p_{0} \int_{0}^{\infty} P(i, t)[1-U(t)] d t+\eta_{i} \quad(i \geqslant 1) \tag{29}
\end{equation*}
$$

for the $T(S V)$-model.

- (22) and (28) we get

$$
\begin{equation*}
p_{i}=\lambda p_{0}\left(1-b_{0}\right)^{-1} \int_{0}^{\infty} P(i, t)[1-U(t)]+\eta_{i} \quad(i \geqslant 1) \tag{30}
\end{equation*}
$$

for the $T(M V)$-model.
When $S(\cdot)$ and $U(\cdot)$ are phase type distributions, explicit results can be obtained for $p_{i}$ $(i \geqslant 1)$ for both models. Note that $p_{0}$ appearing in (29) and (30) is given by

$$
p_{0}=(\mathbb{E}[X])^{-1}(1-\rho)
$$

where (see Loris-Teghem (1990a)) $\mathbb{E}[X]=\lambda \mathbb{E}[U]+b_{0}$ for the former model, while for the latter model $\mathbb{E}[X]=\lambda \mathbb{E}[U]\left(1-b_{0}\right)^{-1}$.

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# THE P-CENTER PROBLEM IN R ${ }^{\text {n }}$ WITH WEIGHTED TCHEBYCHEFF NORMS 

## Blas PELEGRIN PELEGRIN

Dpto. de Matemática Aplicada y Estadística
Universidad de Murcia, 30100-Espinardo. Murcia. Spain

## ABSTRACT

In this paper the p -center problem in $\mathrm{R}^{\mathrm{n}}$ is studied when distances are measured by a weighted Tchebycheff norm. For the 1 -center problem it is proved that the lower bound proposed by Dearing and Francis (1974) is attained and a one step algorithm is given to obtain an optimal solution. Then, an exact algorithm for $\mathrm{p}>1$ is proposed which generalizes the one given by Aneja et al. (1988) for the unweighted rectangular p-center problem on the plane. Finally, a new 2-approximation heuristic polynomial algorithm is given which is $\alpha$ best possible» since for $\delta<2$ the existence of a $\delta$-approximation polynomial algorithm would imply that $\mathrm{P}=\mathrm{NP}$.

1. INTRODUCTION

Let $M=\left\{P_{1}, P_{2}, \ldots, P_{m}\right\}$ be a given set of points in the Euclidean space $\mathbb{R}^{n}$. The p-center problem in $\mathbb{R}^{n}$ seeks p points $C_{1}, C_{2}, \ldots, C_{p}, p<m$, so that the maximal weighted distance between each point $P_{i}$ and its closest point $C_{j}, j=1, \ldots p$, is minimized. The problem is usually formulated as :

$$
\begin{array}{ccc}
\text { (P1) } & \operatorname{Minimize} & \max
\end{array} \quad\left\{w_{i} \quad \min \left\{N\left(P_{i}-X_{j}\right)\right\}\right.
$$

where $N$ denotes a norm function and $w_{i}>0, i=1, \ldots, m$.

Alternatively, the problem can be seen as that of finding a partition $\alpha=\left\{M_{1}, \ldots, M_{p}\right\}$ of the set $M$ into $p$ disjoint subsets, so that the maximum among the maximal weighted distances between the best point coptimal solution to the 1 -center problem with points $P_{i}$ in $M_{j}$ ) and the points in each subset $M_{j}$ is minimized. Then the problem is also formulated as :

| $\begin{aligned} & \text { Minimize } \\ & \alpha \in P(M, p) \end{aligned}$ |  |  |  |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |  |  |  |  |  |  |
| where $P(M, p)$ denotes the set of all partitions of $M$ into |  |  |  |  |  |  |  |  |  |  |  |
| disjoint subsets and $r\left(M_{j}\right)$ denotes the optimal value of the |  |  |  |  |  |  |  |  |  |  |  |
| 1-center problem associated with the points in $M_{j}, j=1, \ldots, p$. |  |  |  |  |  |  |  |  |  |  |  |
| Then the points $C_{1}, C_{2}, \ldots, C_{p}$ are found as optimal solutions to the |  |  |  |  |  |  |  |  |  |  |  |
| 1-center problems associated with the subsets of an optimal |  |  |  |  |  |  |  |  |  |  |  |

The problem arises typically for $n=2$ in the field of Location

Theory when one seeks locations for $p$ facilities, so that the maximal distance or travel time between each demand point $P_{i}$ and its closest facility $C_{j}$ is minimized.Examples of this kind of problem are in the location of emergency services (like fire stations, ambulance and helicopter bases, police stations), location of radio and $T V$ stations, messenger delivery services, etc. The problem can also be found in Cluster Analysis where the aim is to create $p$ groups in a set of objects $M$, so that the maximal dissimilarity between each object $P_{i}$ and the center $C_{j}$ of its group is minimized. The objects are characterized by the values of $n$ variables, and the dissimilarity is measured by a norm function. Examples are found in economics, social sciences, and almost all empirically based disciplines.

For $n=2$, in spite of being $N P$-hard (Megiddo and Supowit'), the problem has recently been solved by exact algorithms for the euclidean norm (Drezner ${ }^{2}, V i j a y^{3}$, Chen and Handler ${ }^{4}$ ), and for the rectangular norm (Watson-Gandy ${ }^{5}$, Drezner ${ }^{6}$, Aneja et al. ${ }^{7}$ ). For arbitrary $n$, due to its complexity, only heuristic algorithms have been used (Drezner ${ }^{2}$, Spath ${ }^{8}$, Eiselt and Charlesworth ${ }^{9}$, Dyer and Frieze ${ }^{10}$, which can be applied for any norm N. Most of these heuristics require the 1 -center problem to be solved many times, which is time consuming for $n>2$. In this case, the centers $C_{1}, \ldots, C_{p}$ are normally constrained to be points in $M$, as happens in other Cluster Analysis problems and in the classical p-median problem in Location Theory.

The aim of this paper is to study the problem when $N$ is given by a weighted Tchebycheff norm, which is defined as:

$$
N(X)=\max \left\{\lambda_{h}\left|x_{h}\right|, h=1, \ldots, n\right\}
$$

where $X=\left(x_{1}, \ldots, x_{n}\right)$ and $\lambda_{h}>0, h=1, \ldots, n$. For $n=2$, these norms are obtained by linear transformation from bidirectional polyhedral norms (Pelegrin ${ }^{11}$ ), and they can then be used in Location Theory when the movement is restricted to two given directions. For instance, taking the transformation $x_{1}=Y_{1}+Y_{2}, x_{2}=Y_{1}-Y_{2}$ it follows that $\left|Y_{1}\right|+\left|Y_{2}\right|=\max .\left\{\left|x_{1}\right|,\left|x_{2}\right|\right\}$, and the Tchebycheff norm is cotained from the rectangular norm. For any $n$, they are obtained by a scaling transformation from the Tchebycheff norm and can be used in Cluster Analysis as criterion distances.

To our knowledge, the problem is studied here for the first time for this type of norm. Firstly, we deal with the 1 -center problem showing that the lower bound given in Dearing and Francis ${ }^{12}$ is reached and solve the problem by a one step algorithm. Then, the algorithm given in Aneja et al. ${ }^{7}$, which solves optimally the problem in $\mathbb{R}^{2}$ with the rectangular norm and $w_{i}=1, i=1, \ldots m$, is generalized to solve the problem in $\mathbb{R}^{n}$ with any weighted Tchebycheff norm and $w_{i}>0, i=1, \ldots, m$ Finally, some heuristic algorithms are considered and a 2-approximation polynomial algorithm is proposed which is "best possible" since, for any $\delta<2$, the existence of a $\delta$-approximation polynomial algorithm would imply that $P=N P$ as shown in Ref. 15 and 16.

The 1-center problem is formulated as follows :

$$
\begin{align*}
& \operatorname{Minimize}  \tag{P3}\\
& X_{X \in \mathbb{R}^{n}} \\
& \quad R(X)= \\
& \operatorname{Max} \quad\left\{w_{i} N\left(P_{i}-X\right)\right\}
\end{align*}
$$

where $X=\left(x_{1}, \ldots, x_{n}\right), P_{i}=\left(a_{1}{ }^{i}, \ldots, a_{n}^{i}\right)$, and $N\left(P_{i}-X\right)=\max \cdot\left\{\lambda_{h}\left|a_{h}^{i}-x_{h}\right|\right.$ $: h=1, \ldots, n\}$. Let $X^{\star}$ and $r^{*}$ denote an optimal solution and the optimal value of (P3) respectively.
$A$ lower bound of $R(X)$ for any norm $N$ is given in Dearing and Francis ${ }^{12}$ by :

$$
B=\max _{i \neq k}\left\{w_{i} w_{k} N\left(P_{i}-P_{k}\right) /\left(w_{i}+w_{k}\right)\right\}
$$

The following property shows that the lower bound $B$ is reached and gives the set $S^{*}$ of optimal solutions to (P3).

Property 1
i) $r^{*}=B$.
ii) $S^{\star}=\left\{X \in \mathbb{R}^{n}: \max \left\{a_{h}{ }^{i}-B / w_{i} \lambda_{h}\right\} \leq x_{h} \leq \min _{1 \leq i \leq m}\left\{a_{h}{ }^{i}+B / w_{i} \lambda_{h}\right\}\right\}$.

Proof:
Let $B(X, r)$ denote the ball centered at $X$ and with radius $r$, i.e., $B(X, r)=\left\{Y \in \mathbb{R}^{n}: \lambda_{h}\left|Y_{h}{ }^{-} X_{h}\right| \leq r, h=1, \ldots, n\right\}$, and let $B_{h}(X, r)$ denote the projection of $B(X, r)$ on the $h-t h$ axis,i.e., $B_{h}(X, r)=$ $\left\{Y_{h} \in \mathbb{R}: \lambda_{h}\left|Y_{h}-X_{h}\right| \leq r\right\}$. Then $B(X, r)=\prod_{h=1}^{n} B_{h}(X, r)$ because $N$ is a weighted Tchebycheff norm (note that this holds for such norms only). Also $r \geq r^{\star}$ iff $\cap\left\{B\left(P_{i}, r / w_{i}\right): P_{i} \in M\right\} \neq \varnothing$.

For $X_{i k}=\left(w_{i} P_{i}+w_{k} P_{k}\right) /\left(w_{i}+w_{k}\right), i \neq k$, it follows that
$B \geq w_{i} w_{k} N\left(P_{i}-P_{k}\right) /\left(w_{i}+w_{k}\right)=w_{i} w_{k}\left(N\left(P_{i}-X_{i k}\right)+N\left(P_{k}-X_{i k}\right)\right) /\left(w_{i}+w_{k}\right)=$

$$
=w_{i} N\left(P_{i}-X_{i k}\right)=w_{k} N\left(P_{k}-X_{i k}\right)
$$

therefore $X_{i k}=\left(x_{1}{ }^{i k}, \ldots, x_{n}^{i k}\right) \in B\left(P_{i}, B / w_{i}\right) \cap B\left(P_{k}, B / w_{k}\right) \quad$ and $x_{h}{ }^{i k} \in B_{h}\left(P_{i}, B / w_{i}\right) \cap B_{h}\left(P_{k}, B / w_{k}\right)$ for $h=1, \ldots, n$. Hence on each $h$-axis $B_{h}\left(P_{i}, B / w_{i}\right) \cap B_{h}\left(P_{k}, B / w_{k}\right) \neq \varnothing$ for any $P_{i}, P_{k} \in M, i \neq k$. From the Helly property ( Rockafellar ${ }^{13}$ ), it follows that $\cap\left\{B_{h}\left(P_{i}, B / w_{i}\right)\right.$ : $\left.\mathrm{P}_{\mathrm{i}} \in \mathrm{M}\right\} \neq \varnothing$.

For $h=1, \ldots, n$, let $x_{h}{ }^{\prime}$ be in $\cap\left\{B_{h}\left(P_{i}, B / w_{i}\right): P_{i} \in M\right\}$, then $X^{\prime}=\left(x_{1}, \ldots, x_{n}^{\prime}\right) \in \cap\left\{B\left(P_{i}, B / w_{i}\right): P_{i} \in M\right\}$. This implies that $R\left(X^{\prime}\right)=B$. Consequently, $r^{\star}=B$ and $S^{\star}=\left\{X \in \mathbb{R}^{n}: R(X)=B\right\}$.

As $X \in S^{*}$ if and only if $w_{i} \max .\left\{\lambda_{h}\left|a_{h}{ }^{i}-x_{h}\right|: h=1, \ldots, n\right\} \leq B$ for all $P_{i} \in M$, the proof is complete.

The lower bound $B$ is also reached in $\mathbb{R}^{2}$ for bidirectional polyhedral norms, however this result is false for the $L_{1}$ norm in $\mathbb{R}^{n}, n>2$, as it happens for $P_{1}=(0,0,0), P_{2}=(1,1,0)$, $P_{3}=(1,0,1), P_{4}=(0,1,1)$ and $w_{i}=1, i=1, \ldots, 4$.

As a consequence of property 1 , the following one step algorithm gives an optimal solution to (P3).

ALGORITHM 1

- Calculate :
$B=\max \left\{w_{i} w_{k} /\left(w_{i}+w_{k}\right) \quad \max \left\{\lambda_{h}\left|a_{h}^{i}-a_{h}{ }^{k}\right|\right\}\right\}$
$i \neq k \quad 1 \leq h \leq n$
$\alpha_{h}=\max \left\{a_{h}{ }^{i}-B / w_{i} \lambda_{h}\right\}$ for $h=1, \ldots, n$
$1 \leq i \leq m$
$\beta_{h}=\min \left\{a_{h}{ }^{i}+B / w_{i} \lambda_{h}\right\}$ for $h=1, \ldots n$ $1 \leq i \leq m$

```
    Choose }\mp@subsup{x}{h}{*}\mp@subsup{}{}{\star}\in[\mp@subsup{\alpha}{h}{\prime},\mp@subsup{\beta}{h}{}], e.g. {\mp@subsup{x}{h}{}\mp@subsup{}{}{\star}=(\mp@subsup{\alpha}{h}{}+\mp@subsup{\beta}{h}{})/2 for h=1,.,n
- Set }\mp@subsup{X}{}{*}=(\mp@subsup{x}{1}{}\mp@subsup{}{}{*},\ldots,\mp@subsup{x}{n}{*})\mathrm{ and }\mp@subsup{r}{}{*}= B . Output \mp@subsup{X}{}{*}\mathrm{ and }\mp@subsup{r}{}{*}\mathrm{ . Stop.
3. AN EXACT ALGORITHM FOR P>1
    In the following, the algorithm given in Aneja et al. is
generalized to solve the problem (P2). Let r* denote the optimal
value of (P2) and rik = w w wh}N(\mp@subsup{P}{i}{}-\mp@subsup{P}{k}{\prime})/(\mp@subsup{w}{i}{}+\mp@subsup{w}{k}{})\mathrm{ for i*k . For any
partition \alpha=(M, ,.M}\mp@subsup{M}{p}{})\mathrm{ , from property 1 it follows that :
                        r(M}\mp@subsup{M}{j}{})=\operatorname{max}{\mp@subsup{r}{ik}{}:\mp@subsup{P}{i}{},\mp@subsup{P}{k}{\prime}\in\mp@subsup{M}{j}{}
                                    i\not=k
    Then, it is not necessary to solve any optimization problem to
evaluate r r , as happens with other norms, and the following
property satisfies.
    Property 2
r** R={ rik}:\mp@subsup{|}{ik}{*},\mp@subsup{P}{k}{*}\inM,i\not=k}
    Let r\inR, and define the graph G(r)=(V,E(r)) as follows:
    V}={\mp@subsup{P}{1}{},\ldots,\mp@subsup{P}{m}{\prime}
    E(r)={( P
    Property 3
    (P2) is equivalent to the problem of finding the minimum value
of r in R , such that }V\mathrm{ , is covered by a set of p cliques (maximal
complete subgraphs) of G(r).
    Proof:
```



```
for }\mp@subsup{P}{i}{},\mp@subsup{P}{k}{}\in\mp@subsup{M}{j}{\prime},j=1,\ldots,p. Therefore, M M,j=1,\ldots,p, are complete
subgraphs of G(r) from which a set of p cliques covering v can
be generated adding (if necessary) points to each subset M M
Conversely, if }\mp@subsup{V}{1}{},\ldots,\mp@subsup{V}{p}{}\mathrm{ are cliques of }G(r) covering V, the
a partition \alpha can be generated,eliminating (if necessary) some
points in }\mp@subsup{V}{1}{}\ldots..\mp@subsup{V}{p}{}\mathrm{ , which satisfy }\mp@subsup{r}{\alpha}{}\leqr . Note that since thi
transformation of partitions into sets of covering cliques may be
done in polynomial time, the problems are polynomially equivalent.
Let \mp@subsup{v}{1}{}\ldots,.\mp@subsup{V}{q}{}\mathrm{ be all the distinct cliques of G(r). There exist p}}\begin{array}{l}{\mathrm{ cliques of G(r) which cover V if and only if the optimal value of}}\\{\mathrm{ the following set covering problem, SCP(r), is less than or }}\\{\mathrm{ equal to p :}}
SCP(r) Minimize z = \Sigma j m
    s.t. }\quad\mp@subsup{\sum}{j}{}\mp@subsup{a}{ij}{}\mp@subsup{x}{j}{}\geq1,i=1,\ldots,
    x
```

where $x_{j}=1$ if clique $V_{j}$ is chosen, 0 otherwise; $a_{i j}=1$ if $P_{i} \in V_{j}$,
0 otherwise. Then, the following algorithm solves (P2) optimally .
ALGORITHM 2
Step 1 : Arrange all the distinct $r_{i k}$ values, $i \neq k$, into an
increasing sequence $: r_{1}<r_{2}<\ldots<r_{t}$.
Step 2 : By a binary search on the list $R=\left\{r_{p}, r_{2}, \ldots, r_{t}\right\}$ find
the smallest $r$ in the list, i.e. $\mathbf{r}^{*}$, for which the optimal
value of $\operatorname{SCP}(r)$ is $p$.

Step 3 : Output $r^{\star}$ and the cliques $v_{j}{ }^{\star}$ corresponding to $x_{j}=1$ in the optimal solution to $\operatorname{SCP}\left(r^{\star}\right)$.


## 4. HEURISTIC ALGOR\|THMS

Most heuristic methods for other location-allocation problems
can be modified to obtain approximate solutions to (P2) for any
norm N. For instance, the "alternate location- allocation
method" (Cooper ${ }^{13}$, and the "exchange method" (Spath ${ }^{14}$ ),
both suggested for the p-median problem, can be used; the
difference is that, due to the different objective function, it is
necessary to use a 1-center algorithm instead of a 1-median
algorithm as a subroutine. These types of method become more

```
efficient for a weighted Tchebycheff norm since it is enough to
use the evaluation of r(Mg) given in (1) as a subroutine
during the iterative procedure instead of a 1-center algorithm,
which is used only at the end to obtain the centers (Algorithm 1
can be used to obtain the centers from the generated partition ).
```

Moreover, for (P2) some $\delta$-approximation polynomial algorithms have been given, which generate a partition $\alpha$ such that $r_{\alpha} \leq \delta r^{*}$ for some given value of $\delta$. For instance, Dyer and Frieze 10 describe a simple heuristic for the p-center problem with an arbitrary metric which is a $\delta$-approximation for $\delta=\min \{3$, $w$ \}, where $w$ is the maximum ratio between the weights of the points in $M$, and Plesnik 15 gives a 2 -approximation polynomial algorithm ("best possible") for the p-center in graphs; both can be used to obtain approximate partitions of (P2) for any norm $N$ from which the centers can be obtained by Algorithm 1.

We propose a new heuristic algorithm to generate an approximate partition, based on properties 1 and 2 , as follows :

## ALGORITHM 3

Step 1 : - Generate any partition $\alpha_{0}=\left(M_{1}{ }^{0}, \ldots, M_{p}{ }^{0}\right)$.

- Calculate $r_{\alpha_{0}}$.
- Set $R=\left\{r_{i k}: r_{i k} \leq r_{\alpha_{0}}, i \neq k\right\}$.

```
    - Make a list L arranging all the distinct rik values
    of }R\mathrm{ into an increasing sequence.
    - Set }\mp@subsup{r}{0}{}=\mp@subsup{r}{\mp@subsup{\alpha}{0}{}}{}\mathrm{ and }\alpha=\mp@subsup{\alpha}{0}{}
Step 2 : If |L|>1 take r as the median value of L. Make
    unlabelled all points in M . Go to step 3.
    Else L = { ro }, STOP. Output rom and \alpha .
Step 3 : - Choose an unlabelled point P }\mp@subsup{f}{t}{}\mathrm{ of maximum weight and set
    M}\mp@subsup{t}{}{\prime}={\mp@subsup{P}{i}{}:\mp@subsup{r}{it}{\prime}\leqr}. Label all the points in M M'.
    - If all points in M are labelled go to step 4. Else go
    to step 3.
Step 4: - Set \mp@subsup{\alpha}{}{\prime}={\mp@subsup{M}{t}{\prime}:\mp@subsup{M}{t}{\prime}}\mathrm{ is generated in step 3}.
    - If |\alpha'| < p set L & { r ikf L : rik \leq r }, \alpha=\alpha' and
    rom}=r.Else set L&{ rik\inL: r& rik }
    -Go to step 2.
    The complexity of step 1 is O(m}\mp@subsup{}{}{2}\operatorname{log}m). Step 2 to step 4 is a
binary search in the set L that finds the minimum value in L
for which the partition (\mp@subsup{\alpha}{}{\prime}}\mathrm{ generated in step 3 satisfies | |'| s p.
As step 3 is O(m
complexity of the algorithm is }O(\mp@subsup{m}{}{2}\operatorname{log m})
```


## Property 4

```
The above algorithm is a 2 -approximation polynomial algorithm for (P2), and gives a lower bound \(r_{0}\) of \(r^{*}\).
Proof :
First, it will be shown that \(r\left(M_{t}{ }^{\prime}\right) \leq 2 r\) for each set \(M_{t}{ }^{\prime}\) generated in step 3 and any \(r \in L\). From the definition of \(M_{t}^{\prime}\) and
```

(1), this is true if $\left|M_{t}{ }^{\prime}\right|=1$ or $\left|M_{t}{ }^{\prime}\right|=2$. Otherwise, let $P_{i}, P_{k}$ $\in M_{t}{ }^{\prime}, i, k \neq t$, as $r_{i t} \leq r, r_{k t} \leq r, w_{i} \leq w_{t}$ and $w_{k} \leq w_{t}$, then
$r_{i k}=w_{i} w_{k} /\left(w_{i}+w_{k}\right) N\left(P_{i}-P_{k}\right) \leq$
$w_{i} w_{k} /\left(w_{i}+w_{k}\right) \quad\left(N\left(P_{i}-P_{t}\right)+N\left(P_{k}-P_{t}\right)\right) \leq$
$w_{i} w_{k} /\left(w_{i}+w_{k}\right)\left(\left(w_{i}+w_{t}\right) / w_{i} w_{t}+\left(w_{k}+w_{t}\right) / w_{k} w_{t}\right) r \leq$
$\left(1+2 w_{i} w_{k} /\left(w_{i}+w_{k}\right) w_{t}\right) r \leq$
$\left(1+\max \left\{w_{i}, w_{k}\right\} / w_{t}\right) r \leq 2 r$,
therefore from (1) it follows that $r\left(M_{t}{ }^{\prime}\right) \leq 2 r$.
Let $r \geq r^{*}$, that is there is a partition $\alpha=\left(M_{1}, \ldots, M_{p}\right)$ with $r_{\alpha} \leq r$. Let the point $P_{t}$ in an iteration of step 3 be in $M_{j}$; as $r_{i t} \leq r$ for any $P_{i} \in M_{j}$ it follows that $M_{j} \subset M_{j}$. . Then the partition $\alpha^{\prime}$ generated in step 3 satisfies $\left|\alpha^{\prime}\right| \leq p$. Consequently, if the partition $\alpha^{\prime}$ generated in step 3 satisfies $\left|\alpha^{\prime}\right|>p$ then $r$ < $r^{*}$. Therefore min $\left\{r_{i k}: r_{i k} \in L\right\}$ is a lower bound of $r^{*}$ in each iteration of step 4 .

At the end $L=\left\{r_{0}\right\}$ and $|\alpha| \leq p$, as then $r_{o} \leq r_{\alpha} \leq 2 r_{0} \leq 2 r^{*}$ it follows that Algorithm 3 is a 2-approximation polynomial algorithm and $r_{0}$ is a lower bound of $r^{*}$.

The proposed algorithm is then "best possible" since for any $\delta<2$ it is known that the existence of a $\delta$-approximation polynomial algorithm would imply that $P=N P$ (Hsu and Nemhauser ${ }^{15}$, Hochbaum and shmoys ${ }^{16}$ ). It generates a partition $\alpha$, as happens with the other algorithms mentioned, so that the centers can be obtained from the generated partition by Algorithm 1 . Also, an upper bound on the related error ( $\left.r_{\alpha^{-}} r^{*}\right) / r^{\star}$ can be obtained as $\left(r_{\alpha} \cdot r_{0}\right) / r_{0}$.

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# A VARIATION OF GRAHAM's LPT-ALGORITHM 

Hannes HASSLER<br>TU Graz<br>Institut für Angewandte Informatik und Kommunikationstechnologie<br>Klosterwiesgasse 32, A-8010 Graz, Austria<br>and<br>Gerhard J. WOEGINGER<br>TU Graz<br>Institut für Mathematik B<br>Kopernikusgasse 24, A-8010 Graz, Austria

## ABSTRACT

We consider the problem of scheduling a set of $n$ jobs nonpreemptively on $m$ identical machines. The goal to minimize $C^{\star}$, the maximum completion time over all jobs. This problem is known to be NP-complete. We present a class of new approximation algorithms that have low running time $O(n \log m$ ) and guarantee the worst case performance of Graham's famous LPT-Algorithm.

## 1 Introduction

We consider the problem of scheduling a set of $n$ jobs nonpreemptively on $m$ identical machines. The goal is to minimize the maximum completion time over all jobs. We denote the maximum completion time in an optimal schedule by $C^{*}$. In 1969, Graham [2] suggested two simple heuristics for this problem. List Scheduling, the first heuristic, lists the jobs in any order and then assigns them in this order to the machines as they become free. The Longest Processing Time algorithm (or LPT, for short) first sorts the jobs by decreasing processing time and then applies list scheduling to them. List scheduling gives a schedule with maximum completion time at most $2-\frac{1}{m}$ times $C^{*}$ and LPT gives a schedule with maximum completion time at most $\frac{4^{m}}{3}-\frac{1}{3 m}$ times $C^{*}$. Both bounds are tight. For the running time, it is easy to see that list scheduling takes $O(n \log m)$ time and LPT takes $O(n \log n)$ time. Thus, if we take the number of machines to be constant, the running time of list scheduling outperforms LPT by a $\log (n)$-factor.

In this note, we present an algorithm that has the same worst case performance guarantee as the LPT-algorithm and has the low running time of list scheduling. More precisely, we construct a sequence of algorithms that lie between list scheduling and LPT and determine their exact worst case performance. The behaviour of the algorithms depends on the number $m$ of machines and on some parameter $k, 0 \leq k \leq 2 m$. For $k \geq 2 m$, our algorithm is at most the LPT-factor $\frac{4}{3}-\frac{1}{3 m}$ off the optimum.

## 2 The Algorithm and the performance guarantee

Let $J$ be a set of $n$ jobs that must be scheduled to $m$ machines and let $k$ be some integer, $0 \leq k \leq n$. We may assume $m \leq n$ as otherwise the problem is trivial. Our approximation algorithm $A(m, k)$ proceeds as follows.
(1) Determine the set $J^{k}$ of the $k$ jobs in $J$ with largest processing times. Form a list that contains the jobs in $J^{k}$ in arbitrary order and append the remaining jobs in $J-J^{k}$ to it.
(2) Apply list scheduling to the constructed list.

It is well known that determining the $k^{\text {th }}$-largest processing time can be done in $O(n)$ time (see e.g. [1]). Thus, finding $J^{k}$ takes $O(n)$ time and appending the remaining jobs takes $O(n)$ time, too. Step 2 is standard $O(n \log m)$ list scheduling, and so we get an overall time complexity of $O(n \log m)$ for algorithm $A(m, k)$.

Now let $C_{k}^{H}$ denote the maximum completion time in a schedule constructed by the algorithm $A(m, k)$ and let $C^{*}$ denote the maximum completion time in an optimal schedule. In Section 3, we will prove the following tight worst case bounds.

| Range of $k$ | Worst case Ratio <br> of $C_{k}^{H} / C^{*}$ | Ratio <br> for $m=2$ |
| :---: | :---: | :---: |
| $0 \leq k \leq(m-1) / 2$ | $2-1 /(m-k)$ | $3 / 2$ |
| $(m-1) / 2 \leq k \leq m-1$ | $2-2 /(m+1)$ | $4 / 3$ |
| $m \leq k \leq(4 m-1) / 3$ | $1.5-1 /(4 m-2 k)$ | $5 / 4$ |
| $(4 m-1) / 3 \leq k \leq 2 m-1$ | $1.5-3 /(4 m+2)$ | $6 / 5$ |
| $2 m \leq k$ | $1.33-1 /(3 m)$ | $7 / 6$ |

That means, the worst case ratio remains constant for all $k \geq 2 m$. However, the running time detoriates and so in practice we will not use large values for $k$. Numerical experiments demonstrated that the average performance of algorithm $A(m, k)$ improves with increasing $k$. Figure 1 graphically illustrates the worst case behaviour of $A(m, k)$ for $k \geq 0$.

## 3 The Proofs

Let $m$ denote the number of machines and let $k$ be the parameter of algorithm $A(m, k)$ as defined in the preceding section. As $k$ is fixed we write $C^{H}$ for $C_{k}^{H}$ to simplify notation. By $x$ we denote the job that is treated last by the algorithm.

First we consider the cases $0 \leq k \leq m-1$. The jobs in $J^{k}$ all have length at least $x$. W.l.o.g. the jobs in $J^{k}$ are scheduled to the last $k$ machines. Consequently, at the time before job $x$ is scheduled by the algorithm, the total load of the first


Figure 1: The worst case performance of algorithm $A(m, k)$
$m-k$ machines is at most $m C^{*}-(k+1) x$. Job $x$ is scheduled to the machine with smallest load at this time and the smallest load is less or equal the average load of the first $m-k$ machines. This gives

$$
\begin{equation*}
C^{H} \leq\left(m C^{*}-(k+1) x\right) /(m-k)+x=\left(m C^{*}+(m-2 k-1) x\right) /(m-k) \tag{1}
\end{equation*}
$$

For $0 \leq k \leq(m-1) / 2$, the coefficient of $x$ in inequality (1) is nonnegative. Using $x \leq C^{*}$, we get that $C^{H} / C^{*} \leq 2-1 /(m-k)$ holds, the claimed result for $0 \leq k \leq(m-1) / 2$. For $(m-1) / 2 \leq k \leq m-1$, the coefficient of $x$ in inequality (1) is nonpositive. In this case for $x \geq m C^{*} /(m+1)$, inequality (1) immediately implies $C^{H} / C^{*} \leq 2-2 /(m+1)$. For $x \leq m C^{-} /(m+1)$, we use another averaging argument. At the time before job $x$ is scheduled, the total load over all machines is less or equal $m C^{*}-x$. Hence

$$
\begin{equation*}
C^{H} \leq\left(m C^{*}-x\right) / m+x=C^{*}+(m-1) x / m \leq(2-2 /(m+1)) C^{*} \tag{2}
\end{equation*}
$$

holds, and we have finished the proof for cases $0 \leq k \leq m-1$.
Next, we treat the cases $m \leq k \leq 2 m-1$. We define $k^{\prime}=k-m \geq 0$. First we observe that

$$
\begin{equation*}
x \leq C^{*} / 2 \tag{3}
\end{equation*}
$$

must hold. If $x$ is greater than $C^{*} / 2$, there are $k+1>m$ big jobs of length greater $C^{*} / 2$. Then in the optimum schedule there exists a machine with two big jobs and load greater $C^{*}$, a contradiction. After the algorithm scheduled the $m+k^{\prime}$ jobs in $J^{k}$, there are at least $m-k^{\prime}$ machines that received exactly one job of $J^{k}$ (Every machine gets at least one job. If there are only $m-k^{\prime}-1$ machines with one job, the remaining $k^{\prime}+1$ have at least two jobs each, and we get a total job number of at least $m+k^{\prime}+1>k$ jobs in $J^{k}$, a contradiction). We rearrange the sequence of machines such that the $m-k^{\prime}$ machines with one job of $J^{k}$ range first. The last $k^{\prime}$ machines process the remaining $2 k^{\prime}$ jobs. Consequently, at the time before job $x$ is scheduled by the algorithm, the total load of the first $m-k^{\prime}$ machines is at most $m C^{-}-\left(2 k^{\prime}+1\right) x$ and

$$
\begin{equation*}
C^{H} \leq\left(m C^{-}-\left(2 k^{\prime}+1\right) x\right) /\left(m-k^{\prime}\right)+x=\left(m C^{-\prime}+\left(m-3 k^{\prime}-1\right) x\right) /\left(m-k^{\prime}\right) \tag{4}
\end{equation*}
$$

holds. For $0 \leq k^{\prime} \leq(m-1) / 3$, the coefficient of $x$ in the right hand side of inequality (4) is nonnegative. Using (3), we derive $C^{H} / C^{*} \leq 1.5-1 /\left(2 m-2 k^{\prime}\right)$ and we are finished. For $(m-1) / 3 \leq k^{\prime} \leq m-1$, the coefficient of $x$ in the right hand side of inequality (4) is nonpositive. We distinguish between the two subcases $x \leq m C^{*} /(2 m+1)$ and $x \geq m C^{*} /(2 m+1)$. In the case $x \geq m C^{*} /(2 m+1)$, we can use inequality (4) again and we get $C^{H} / C^{*} \leq 1.5-3 /(4 m+2)$. In the case
$x \leq m C * /(2 m+1)$, we consider the time before job $x$ is scheduled. The total load over all machines is less or equal $m C^{*}-x$. We derive that

$$
\begin{equation*}
C^{H} \leq\left(m C^{*}-x\right) / m+x=C^{*}+(m-1) x / m \leq(1.5-3 /(4 m+2)) C^{*} \tag{5}
\end{equation*}
$$

holds, and the cases $m \leq k \leq 2 m-1$ are finished, too.
Finally, we assume that $2 m \leq k$ holds. It is easy to see that $x \leq C^{*} / 3$ must hold. Otherwise, there are $k+1 \geq 2 m+1$ jobs of length $>C^{*} / 3$. In the optimum schedule, at least one machine must contain at least three of these jobs; this machine would have load $>C^{*}$. Therefore, we may reuse inequality (5) and plugging in $x \leq C^{*} / 3$ instead of $x \leq m C^{*} /(2 m+1)$, we derive $C^{H} / C^{*} \leq$ $4 / 3-1 / 3 m$. This completes the proof of all worst case bounds claimed in Section 2. $\square$

To prove the tightness of our bounds, we exhibit the following five sets of examples. In each example sequence, the first $k$ elements are the largest $k$ elements in the sequence. Hence, algorithm $A(m, k)$ may skip step (1) and use the sequence exactly in the given ordering.

Example 1. $(0 \leq k \leq(m-1) / 2)$ First there are $k$ jobs of length 1 , then ( $m-k)(m-k-1)$ small jobs of length $1 /(m-k)$, and finally job $x$ of length 1 . By scheduling all small jobs on $m-k-1$ machines, we see that $C^{*}=1$ holds. Our approximation algorithm equally schedules the small jobs on $m-k$ machines, and finally puts $x$ on one of them. Hence, $C^{H}=2-1 /(m-k)$.
$\square((m-1) / 2 \leq k \leq m-1)$ Our list consists of $m-1$ big jobs of length $m /(m+1), m$ small jobs of length $1 /(m+1)$ and the big job $x$ of length $m /(m+1)$. The optimum solution simply matches every big job with a small job and gives $C^{*}=1$. Our algorithm schedules all small jobs to one machine. Then job $x$ makes $C^{H}=2 m /(m+1)$.
( $m \leq k \leq(m-1) / 3$ ) We use $k$ big jobs of length $1 / 2$, then $(k-m)(k-m-1)$ small jobs of length $1 /(2 k-2 m)$, and finally the big job $x$ of length $1 / 2$. Similarly as in example 1 , we see that $C^{*}=1$ and that $C^{H}=1.5-1 /(4 m-2 k)$. $\square$
$((m-1) / 3 \leq k \leq 2 m-1)$ The first $2 m-1$ jobs are big jobs of length $m /(2 m+1)$, then there are $m$ small jobs of length $1 /(2 m+1)$ and job $x$ of length $m /(2 m+1)$. The optimum schedule assigns to each machine two big and one small job, the approximation schedule assigns all small jobs to one machine. Hence, $C^{*}=1$ and $C^{H}=1.5-3 /(4 m+2)$.
( $2 m \leq k$ ) Here we use Graham's lower bound example [2] for the

LPT-algorithm. We have jobs $p_{1} \ldots p_{2 m}$ such that $p_{i}$ has length $2 m-\lceil i / 2\rceil$ and $p_{2 m+1}=x$ has length $m$. It is easy to check $C^{*}=3 m$ and $C^{H}=4 m-1$. $\square$

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# MODELLING CORPORATE VULNERABILITY: YET ANOTHER EMPIRICAL ATTEMPT 

## D.J.E. BAESTAENS

Erasmus University Rotterdam
Dept of Finance
Burg. Oudlaan 50, 3062 PA Rotterdam, The Netherlands

## ABSTRACT

This paper argues that potential weaknesses of MDA could be avoided by using the Mahalanobis Distance ( $d^{2}$ ). The model's uniqueness allows for a more detailed analysis of corporate failure triggers (conversely shock absorbers) at the corporate level. Incidentally, the main input for the distance analysis, the sampie's dispersion matrix, could be used to quantify the distance's sensitivity to changes in the (co)variances of so-called policy variables Such analysis is likely to benefit central bank regulators (monetary \& fiscal policy) as well as corporate management (interest rate sensitivity).

## 1. Introduction

Company failure is an emotionally event, unfortunately not really accounted for by a solid theoretical framework ${ }^{2}$. As Barnes (1987) pointed out, researchers may well be testing the hypothesis that the event may be explained by the statistical behaviour of ad hoc generated datasets.

The first attempts to quantify the bankruptcy process in the US were presented in the seminal Altman (1968,1977,1984) papers. Taffler (1977,1980,1982) was among the first to introduce and improve upon the Altman MDA-method ${ }^{3}$ in the UK corporate environment. These original contributions have since been the subject of elaborate testing procedures and of academic debate (see eg. Ezzamel and Mar-Molinero,1990; Karels and Prakash,1987) whereby the discussion focuses on the misspecification of the MDA model. Are researchers justified in applying a multivariate statistical technique while disregarding its restrictive assumptions ?
In theory, the use of MDA requires the prevalence of at least two easily identifiable and mutually exclusive groups of observations, multivariate normality of the cases under study, specification of the a priori probabilities of occurrence of each group and dispersion matrices' equality in case of linear discriminant analysis (LDA).

In this paper we assert that the first two requirements are very restrictive indeed. Because of these constraints the MDA method may not be the most appropriate method to describe and predict failure. Consequently, we will suggest a potential alternative to MDA. Finally, we will apply the alternative to a sample of UK firms from the Printing \& Publishing Industry with a view to identifying entities deviating from the industry average.

## 2. Limitations of MDA

## 1. Group Composition

Although the first requirement does not usually constitute an obstacle, identifying distinct groups of failed and nonfailed firms may be more treacherous than it seems. How does one define failure when the failure process is conceptually nebulous ? A majority of studies appears to adopt some ex post empirical record as failure criterion, be it renegotiation with creditors, bond default, stock exchange delis-

[^0]ting or decision to file for reorganisation under Chapter xI of the Federal Bankruptcy Act.
While analysis' of the real world variable clearly indicates the category the firm belongs to, classification of the cases becomes a biased process as the failure category becomes unnecessarily stratified. To exemplify, we have not come across many MDA studies investigating the question why some firms are more successful than others in a given setting. Metaphorically, analysing failures may be like investigating the Black Death. What matters is not the detection of the pathogen by carrying out post-mortem analyses but the identification of the survivors' defence mechanisms. We fear that restraining the failure class sample to those entities that actually failed (whatever the failure definition) may lead to the repudiation of the evidence provided by those firms that were infected with the failure-disease but managed to get rid of the malady.
We therefore prefer an a priori undefined approach to the ex post empirical classification routine. The advantages of the a priori approach are fourfold. First, we do not have to hunt for a specific, necessarily subjective (in the absence of a failure theory) failure criterion. Second, by pooling all firms we avoid the obligation of having to construct a non failed sample consisting of really healthy firms unlike Taffler (1982). Third, no tests are needed to verify the equality of the variance-covariance matrices across groups. Finally, our procedure possesses the implied benefit of not having to deal with the issue of determining a priori probabilities of the event occurrence.

## 2. Multivariate Normality

The normality condition sometimes appears to be interpreted as imperative (Eisenbeis,1977; Karels and Prakash,1987; Taffler,1982) and sometimes as accessory (Barnes,1982; Richardson and Davidson,1983). This ambiguity may stem from the lack of a standard test for multivariate normality compounded by the fact that univariate normality does not guarantee multivariate normality. As non normality may be caused by the presence of outliers, the conventional statistical approach has always been to identify such outliers with a view to deleting or separating them from the rest of the data under study (Anscombe,1960; Collett and Lewis,1976; Hartwig and Dearing,1979).

Our present approach on the contrary may be somewhat novel in the sense that outliers are treated as the major information source (Ezzamel and Mar-Molinero,1990; Howell,1989).

4 Sometimes euphemistically called Winsorization (changing an outlier's value to that of the closest non outlier) or Trimming (removal from the sample an equal number of the smallest and largest observations).

We assume that events in which either individual variables or the relationships between variables take unusual values or are distorted into unusual levels (as measured by statistical criteria) may reflect strain in or upon the economy, respectively the individual firm. This strain may or may not endure after the stress is removed ${ }^{\text {s }}$. Here we are interested in extreme states in their own right, though we accept that some of them will be due to spurious observations or chance events. Where there are many variables, as here, multivariate methods are well developed only for the normal distribution, as discussed by Bacon-Shone and Fung (1987), and this paper uses the Mahalanobis distance ( $\mathrm{d}^{2}$ ) and Hotellings $\mathrm{T}^{2}$ as representative of such methods. Although $d^{2}$ assumes either a multivariate normal or elliptical distribution (Mitchell and Krzanowski, 1985), we are hoping that our data are not jointly normally distributed since we are actively seeking outliers from such a distribution. In this sense, we believe we are among the first to apply the Mahalanobis distance in its own right to the issue of identifying sick firms.

## 3. $\mathrm{d}^{2} \& \mathrm{~T}^{2}$ : Alternatives to MDA ?

### 3.1. Presentation of $d^{2}$

The joint distribution of normally distributed individual variables is often multivariate normal. Figure 1 shows a resulting ellipse of uniform probability density for two such variables $X$ and $Y$, standardised to equal standard deviations.

A joint confidence region for $X$ and $Y$ is elliptical. This region is not the intersection of the two univariate confidence levels at the same significance level (circle in Figure 1). We call this circle an "uncorrelated" confidence region and points outside it "uncorrelated" outliers. Points in Regions $I$ and II are respectively outliers and inliers for both the ellipse and the circle, whilst points in the Regions III are inliers to the ellipse but outliers to the square. Points in Regions IV are outliers to the ellipse but inliers to the square.

Conventional regression analysis does not yield confidence regions equivalent to the ellipse - projection into the $X$ space makes predictions of $Y$ conditional on $X$ rather than absolute, so that the confidence region is hyperbolic and unbounded before $X$ is observed, and "unusual" states of $X$ are not recognised on observation. In general, the absence of a theoretical framework disallow researchers to give any particular set of $X$ variables special status as independents.
s In engineering disciplines, stress and strain denote distinct realities. Stress refers to the event that causes the distortion to happen while strain is the equivalent of the distortion.

Figure 1: Possible confidence regions using uncorrelated and correlated multivariate criteria.

|  |
| :---: |

To deal with this we note that points on the ellipse in Figure 1 share not only the same probability density, but also the same Mahalanobis Distance from the mean observation. The Mahalanobis Distance of a single multivariate observation from the mean observation of a sample of $n$ observations can be estimated using:
$d^{2}=\sum_{1=1}^{p} \sum_{j=1}^{p}\left(x_{1}-\bar{x}_{1}\right) c^{1 j}\left(x_{j}-\bar{x}_{j}\right)$
where
$\bar{x}_{1}$ is the mean of the $i$ th variable
and $c^{1 j}$ is the element in the ith row and jth column of the inverse of the variance-covariance matrix $C^{-1}$.
$d^{2}$ follows a Chi squared distribution with $p$ degrees of freedom (Manly,1986). Figure 1 suggest that in correlated data sets this test will be the most useful and is likely to detect a set of outliers distinct from uncorrelated outliers. A chosen cut off value of the Mahalanobis Distance $d^{2}$ separates observations between Mahalanobis Outliers (Regions I plus IV in Figure 1) and Mahalanobis Inliers (Regions II plus III in Figure 1). We assume temporarily that a d ${ }^{2}$ value has been chosen so as to give these regions convenient relative probabilities, and interpret them as follows. Points in Regions III are events where at least one variable is outside its individual confidence interval, but the joint value is not. Such events we call "structure preserving" in the sense that the expected correlation structure is pre-
served. In contrast, points in Regions IV we call "structure violating", because although neither variable is outside its individual confidence interval, the expected correlation of the two variables is violated. Events in Region I may or may not violate correlation structures.

### 3.2. Decomposition of the Mahalanobis distance

In a $p$ dimensional observation each of the $p(p-1) / 2$ pairs of variates may show either structure violating or structure preserving behaviour. It is desirable to have a means of inspecting this behaviour directly, and one which does not make use of an arbitrary $d^{2}$ to discriminate between Structure Violation and Structure Preservation.

We can express $d^{2}$ in equation (1) as the sum of all the elements of a (pxp) matrix $R$ where $F_{1 j}=x_{1} x_{1} V_{1,}^{1}$ (Kendall and Stuart, 1983) so that
$d^{2}=\Sigma_{1 j} \quad F_{1 j}$
Diagonal elements of $F$ represent a weighting of a single squared deviation of the ith element (i.e. variable) from its mean, and off-diagonal elements represent a weighting of the product of the deviations of the ith and jth variables from their respective means. Since $F$ is symmetrical we can conveniently combine off diagonal $F_{11}$ and $F_{j 1}$ in a single cell by defining $T$ ( $p \times p$ ) such that
$\begin{array}{ll}\text { for } i=j, & T_{1 j}=F_{1 j}, \\ \text { for } i<j, & \mathbf{T}_{1 j}=2 F_{1 j} \\ \text { for } i>j, & T_{i j}=0\end{array}$
Diagonal elements of $T$ show the contribution to $d^{2}$ of the individual deviation of each variable in isolation, while subdiagonal elements show the specific contribution to $d^{2}$ from each variable's interaction with each single other variable. Diagonal elements in $T$ and $F$ are by definition positive, but off diagonal elements can take either sign.

A negative element $T_{i j}$ for $i<j$ indicates that the joint deviations of $x_{1}$ and $x_{1}$ in this observation are less unlikely than the diagonal elements (their individual unlikelihoods) would suggest, and the fact that they are varying in the expected joint direction is "Structure Preserving". Such terms reflect the fraction of the joint variation of $x_{1}$ and $x_{\text {, }}$ that can be predicted from their correlation, and is akin to the "sum of squares explained" in ANOVA. They may correspond to events in Regions II and III of Figure 1 for the $x_{1}$ and $x_{1}$ concerned. Conversely a positive value for $T_{i 1}$ for $i<j$ indicates that an expected positive or negative correlation has been reversed, and this unexpected joint state of $x_{1}$ and
x , we can call "Structure Violating". It can correspond to an event in Regions IV or I (for variables $x_{1}$ and $x_{1}$ only).

A zero value of $T_{11}$ can occur for $i<j$ if $x_{1}$ and $x_{1}$ are uncorrelated in the sample as a whole ( $\mathrm{T}_{11}=\mathrm{F}_{11}=0$ for $i$ not equal to $j$ ), or if standardised $x_{1}$ or $x_{1}$ or both are close to zero in this observation. Such observations are Structurally Neutral, or Uninformative.

### 3.3. Matrix simplification of $F$ and $T$ to $F$. and $T$.

So far, the triangular matrix $T$ represents nothing more than a rearrangement of the full contribution matrix $F$. The value for $d^{2}$ was not at all affected. As matrices $F$ and $T$ in section above contain $p(p-1) / 2$ different entries per observation, we attempted to simplify these matrices by setting to zero all elements in $F(T)$ whose absolute value was smaller than a filter value s. The resulting matrix can be called $\mathrm{F}_{\mathrm{s}}$ ( $\mathrm{T}_{\mathrm{s}}$ ) , and the approximated Mahalanobis Distance is $\mathrm{d}^{2}$., where
$d^{2} \approx d^{2}=\sum_{1,1-1}^{p} F_{(0,1)}$
and
$d^{2} \approx d_{s}=\sum_{1,1=1}^{p} T_{(s)+1}$

We used no theory to set $s$ but investigated the sensitivity of $d^{2}$, to $s$, and selected the largest value of $s$ for which $d^{2}$. seemed a close and stable approximation to $d^{2}$. This simplified interpretation as well as providing a heuristic to avoid over interpreting the large noise content of a $p$ dimensional observation.

### 3.4. Classification of entries in E .

Given the simplest acceptable $F_{\text {, }}$ (or $\mathrm{T}_{\mathrm{s}}$ ), we sorted the variables in descending order of the joint net contributions to $d^{2}$. (that is, the column, respectively row, totals of $F$ ). This partitioned $F_{\text {. }}$ into three regions: Block $A$, where all totals were positive, Block B where the sums were zero and Block $C$ where all aggregates were negative. We call all variables with nonzero diagonal entries Main Variables, since they matter in their own right.
Block A variables contribute positively to $d^{2}$ and are therefore acting in a Structure violating way. These variables can be observed to reinforce $d^{2}$ in their own right (diagonal entry larger than zero resulting in the classification of this variable as Main Variable) and/or in interaction with other variables (off-diagonal entries larger than zero). The contribution of Block B variables remains neutral or uninformative as all diagonal and off-diagonal elements are set to zero. While the net contribution by Block $\mathbf{C}$ variables
can be classified as structure preserving (i.e. distance reducing), their individual contributions are never negative (diagonal elements are weighted squared deviations from respective means) implying that the variable interactions must more than offset these individual excursions, Again a nonzero positive deviation results in a classification as a Main Variable.

### 3.5. Hotelling's $T^{2}$ to search for longer lasting outlier episodes

Mahalanobis outliers for large enough $d^{2}$ are rare in the sample. They may at times be caused by chance events, or invalid measurements, or valid measurements of relationships that have no substantive economic importance. However they may also reflect important structural effects. For example the observed overall correlation may be due to forces that tend to suppress certain joint states of the variables.
If the system is driven suddenly into a "non favoured" state by some shock or stress, and the corrective forces do not take full effect within a single sampling interval lagged effects may occur.

If such lagged effects are present, some Mahalanobis outliers will form part of extended episodes, in which extreme states in $p$ space are only gradually approached and/or gradually departed from. While the squared Mahalanobis Distance was convenient for testing individual outlier observations, we used the Hotelling's $\mathrm{T}^{2}$ to test whether the mean of a group of $p$ dimensional observations differs from the mean of $a$ second group (the remaining observations) assumed to have the same covariance matrix. For a null hypothesis of equal sample means the relevant test statistic reduces to an $F$ distribution with $p$ and ( $n-p-1$ ) degrees of freedom, in the notation of equation (5) below (Manley,1986; Krzanowski,1988).
$T^{2}=n_{1} n_{2} /\left(n_{1}+n_{2}\right) \sum_{1=1}^{p} \sum_{j=1}^{p}\left(\bar{x}_{11}-\bar{x}_{21}\right) c^{11}\left(\bar{x}_{1 j}-\bar{x}_{2 j}\right)$

## where

$\overline{\mathbf{x}}_{11}=$ the mean of the ith variable for group 1 $p=$ number of variables
$n=$ pooled sample size or ( $n_{1}+n_{2}$ ), and
$c^{1}=$ the element in the ith row and jth column of the inverse of the pooled within group covariance matrix

In order to avoid bias by each outlier itself, we omitted it from each putative episode. We searched by trial and error for subsets of the total sample, contiguous with but not including the outlier(s), which differed significantly from the total sample.

Clearly there is some risk of "Data Mining" in such a search, and we have not attempted to derive an exact correction for it. Explicitly dynamic methods were not used because of the scarcity of degrees of freedom, and because we did not wish to assume uniform simple dynamics throughout the sample.
4. Sumary of hypotheses

1. Extreme values or Outliers will be found in economic data sets which are wholly or mainly multivariate normal in their distribution.
2. Mahalanobis Distances of outliers (or inliers) can be decomposed to show the contributions made by each individual variable and by each pair of variables. The latter can in turn be divided into "Structure Preserving", "Structure Violating" and "Structurally Neutral" behaviour for this pair of variables. No specific predictions are made, but the structure of actual outliers are assumed to be of interest to economic modellers .
3. Mahalanobis Outliers will sometimes be a part of longer Dynamic episodes, in which either the approach to or the retreat from an outlier value in Mahalanobis space spreads over several sampling intervals.
4. Multivariate distributions of data on an economy will not be static, but will show, in addition to outliers, signs of sustained structural change (perhaps corresponding to intuitively identifiable periods of economic history). Such secular changes will increase the scatter of a fitted static distribution, and so reduce the power of tests for outliers, but Mahalanobis extreme values will still occur, and outliers strong enough to be detectable in these conditions may have substantive meaning.

## 5. Findings

### 5.1. Variable and Industry Selection

We selected a sample of 13 companies listed on the LSE and for which the FT provides daily financial coverage under the heading Newspapers \& Publishers. Our sample constitutes about 34\% of the total industry with strong emphasis (about 60\%) on the newspapers \& periodicals segment. The companies are given in table 1.

Table 1: Selected Companies and Code

|  | Black (A.f C. | blac |
| :---: | :---: | :---: |
| 2 | Maxwoll Cowna. Corp. | Muxc |
|  | Trinlity int mid. | TRIN |
|  | utd. Nowspapera | urws |
| 6 | Hawe Int spac.DIv. | MEws |
|  | Porter | PSUN |
| 8 | Read Intarnational | RezD |
|  | Haynor | HYNES |
| 10 | Brintol Eve.poe | BRTL |
| 12 | erap | empr |
| ${ }_{13}^{12}$ | Nows Corp | meusc |

Annual data from 1984 to 1989 (and where possible 1990) were compiled from Datastream. The Printing \& Publishing industry was selected because of its alleged homogeneity in terms of cost structure and of the press coverage some industry members were/are receiving. The small number of companies (in absolute sense) excluded every attempt to restrict the sample to only those companies with the same year end. We agree with Gonedes (1973) that such restriction may be viewed as a sample stratification in the sense that the selected companies may share some characteristics that differ from those companies with different year ends.

Since we are interested in the health (or vulnerability) of individual firms over time relative to the industry considered, we pooled the cross-sectional data on individual firms resulting in 82 observations ${ }^{6}$.

Raw data were collected in seasonally adjusted form and where appropriate differenced to remove time trends. Table 2 lists the variables and codes. It can also be seen that most variables depart from a normal distribution on the Lilliefors test (1967), a variant on the Kolmogorov-Smirnov test when the population parameters are unknown, at the $5 \%$ significance level.

[^1]Table 2: Selected Variables

| No category | variable | Code | K-S test |
| :---: | :---: | :---: | :---: |
| 1 Company specific |  |  |  |
| 1 profitability | operating margin | profmar | Yes |
| ${ }_{3}^{2}$ CApITAL Structure |  | corrat | NO |
| 4 turnover patio | SALEs / NET CUR ASSETS | tormcuas |  |
| 5 Productivity | tax / Pretax profit | taxpat | yes |
| 2 Printing a publishing | madustry related |  |  |
| enrnings index | avg earmitas employees pap industry | EARNRAT | мо |
| trade terks | TO AVG RARAINGG MANUFACTURING IND | tradet | но |
| - productivity index | PPI MATERIALS PURCHASED BY IND output indostry / output mandif ind volume terms | RELOUTP | мо |
| 3 economy mide |  |  |  |
| 9 MARKET CONPIDENCE 10 NATIOHAL OUTPUT | SP / ft all senre index GDP (SA, PERCENTAGE CHANGES) | gPFTA GD | $\underset{\mathbf{y E S}}{\mathbf{N O}}$ |

The bad results on the normality tests, while not very encouraging, must not be dramatised.
First, a non normal distribution represents a well known property of most financial ratios (Karels and Prakash, 1987; Lev and Sunder,1979). Following Barnes (1982) we did therefore not attempt to apply standard Box and Cox transformations for right skewed data using the logarithm, square roots and cubic roots. Negative data values prevented the use of these transformations. Of course we could have added a constant where necessary or we could have applied the reciprocal transformation but we feared that too much data manipulation might have destroyed the ratios' informational content. The univariate non normality, of course, may have an adverse impact on the presence of multivariate outliers implying the need to consider the them with great caution, although univariate normality by itself does not guarantee multivariate normality.
Second, assuming the homogeneity of our sample, it may be argued that for a large sample such as ours a small violation of the normality assumption is sufficient for it to be rejected.

Unlike most MDA studies, we attempted to incorporate three possible determinants of corporate well-being and their interaction effects in one dataset. A firm's financial health can be viewed as a function of firm specific, industry related and/or economy wide elements. Obviously, each of these components could be emulated by numerous variables and ratios. In order to select a meaningful variable set for each component, factor analysis (PCA) may be called for to reduce the number of relevant variables. Here, we selected the variables on the basis of their perceived popularity in the literature.

### 5.2. Hypothesis 1

Table 3 shows the set of five outliers all at $1 \%$ significance. On the null hypothesis one might expect one or two such outliers in a sample of 82 observations.

Table 3: Identified Outlier Set


Koy to coig: dialte rofar to Year


All significant outliers (and for that matter the first ten outliers) refer to either 1984 or 1985. The nonuniform distribution of the outliers over time may indicate that the industry or parts thereof was under severe strain during that period. It would be useful to check the evidence of this assumed strain by examination of the relevant annual reports of the identified firms. Hypothesis 4 (see 18) appears to confirm our expectation of a nonstationary mean for the Mahalanobis model implying a structural break within the economy or industry under investigation. To the extent the break does not alter the estimated dispersion matrix, the model's validity remains safeguarded.

As aforementioned, the outliers may represent real phenomena or invalid measurements. The findings for hypothesis 2 on 16 may suggest measurement problems regarding the highest $d^{2}$ values.

Figure 2 and Figure 3 show all outliers over time ranked by firm. It can be seen that the outlier for firm 7 (PSUN) appears to be part of a longer dynamic episode in which both the approach to and the retreat from the outlier value spreads over several sampling intervals. In contrast, the values for firms 10 (BRTL) and 13 (DMGT) seem to belong to an episode in which the approach to, respectively the retreat from the outlier marks a rather abrupt process.

Figure 2: Distances Firms 1 to 7


Figure 3: Distances Firms 8 to 13


We now decompose the $d^{2}$ values by rearranging the variables in such a way that we gain insight in their individual and joint behaviour.
Since our variable set is not very large ( $p=10$ ), the 45 pairs can be evaluated without necessarily simplifying the contribution matrix (F). Table 4 shows contribution matrices ranked by the net variable contribution to $d^{2}$ and $d^{2}$. (ie. column, respectively row, labelled sum) for firm 7 (PSUN) for 1985.

The $d^{2}$ contribution matrix shows that all variables are active, that is display nonzero entries. This finding which was confirmed by the analysis of all $F$ matrices supports our claim that it is dangerous to rely on univariate criteria when assessing corporate vulnerability. Moreover, the finding demonstrates the weakness of those techniques such as MDA that attempt to identify failure trigger points by the sole use of internal, company specific factors. External (industry and economy) factors are not to be omitted from the analysis.

It can be seen that both $d^{2}$ and $d^{2}$ appear to be driven (Block A) by the solitary contribution of the tax ratio variable and the joint contributions of the GDP variable. Structure preserving behaviour (Block $C$ ) is displayed by the trade terms, turnover to net current assets and liquidity ratios. The remaining variables do not contribute in any way to $d^{2}$ : and are therefore classified as Block B variables.
All Block $A$ variables are Main variables, i.e. important in their own right whereas one Block $C$ variable (Trade terms) figures as a nonmain variable (i.e. unimportant in its own right).

The decomposed contribution matrix $F$ could be used as input for the selection of appropriate "policy variables" by regulatory authorities or industry watchdogs. To understand the variables joint behaviour within the system, estimated variance and covariances could be altered to quantify the effects (that is, changes in $d^{2}$ ) of potential policy changes. The aim would then be to minimise $d^{2}$ or to bring it below some agreed reference value. To exemplify, we doubled the variance of the TAXRAT variable mentioned in Table 4 (authorities indirectly control this variable through the tax rate). Table 5 gives the result.

Table 5: Quantification of Intervention Impact PSUN 1985

|  | ${ }_{7 \times 18}$ |  |  |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |  |  |  |  |  |  |
| git |  |  | -0.22 | .oo |  |  |  |  |  |  | 3:49 |
|  | -0.08 | -0.63 | -:30 | -0:32 | -0: | -0.3 | -0: | -0.01 | -12 | -0.1 | 16 |
| TuRN |  |  |  |  |  |  |  |  |  |  |  |
|  |  |  |  |  |  |  |  |  |  |  | 1 |
| ${ }_{\text {ca }}$ |  |  |  |  |  |  |  |  |  | .99 | -0.04 |
| gun | 3.6 |  |  |  |  |  |  |  |  |  |  |

Due to the intervention, this observation has become an inlier. Note that the TAXRAT variable is still the most important Main variable. Obviously, to fully assess the impact of government's intervention, it is necessary to compute all $d^{2}$ values again to ascertain whether or not outliers have become inliers or vice versa. We hope our simplified example has demonstrated the model's potential usefulness in identifying target, intermediary and policy variables without having recourse to monetarist or postkeynesian theories.

Potential measurement problems were signalled above. Analysis of the largest outlier (CODE 585: UNWS 1985, $\mathrm{d}^{2}=79$ ) revealed that almost all of the distance value could be related to the individual contribution of one variable, the borrowing ratio. Table 6 gives the decomposition.

Table 6: Decomposition of the largest outlier, UNWS 1985

|  |  | ${ }_{11} 1^{1}$ |  | Turnclias |  | -raiet | ${ }^{1} 1$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| ciomp |  |  | a.22 0.02 0.02 |  | -0. | O-1.00 -0.05 -0.05 | -0. | -0. |  | -1.24 | 0.17 0.15 |
| turacuas |  |  |  |  | -o. |  |  |  |  |  | 04 |
| ${ }_{\text {Rement }}$ |  |  | -0.00 | -0.01 | -23 | -0.01 |  |  |  | O. 15 | .os |
| ${ }_{\text {TR }}^{\text {TRARPA }}$ | -i. ${ }^{-1} 5$ | -0, | ${ }^{-0.085}$ | ${ }^{-0.02}$ | -0.01 | 号.78 | 0. | - | - 0.22 | -0. ${ }^{21}$ | -0.09 |
| endmrat |  | 15 | -0.07 | -0.01 | -0.10 | - | - 0.15 | 3.36 | - | -0.51 | -0.18 |
| SPFta | -1:73 | - 24 | -0.08 | -0.02 | -0.15 | - | -0.03 | -0.31 | 0.11 | 1.02 | -0.56 |
| un | 0.10 | 0.17 | 0.05 | 0.04 | -0.00 | -0.04 | -0.07 | -0.17 | -0.18 | -0.5 |  |

Since the totals for all variables but the borrowing ratio hardly deviate from zero, we checked the accuracy of the BORRAT variable's value. As the raw value equalled $\mathbf{- 1 2 9 . 4 , ~ w e ~}$ had Datastream confirm this value and its sign ${ }^{7}$.
The $F$ matrix could be simplified such that almost all offdiagonal elements became zeros without impairing the original $d^{2}$ value. The importance of the diagonal entries (BORRAT and EARNRAT) indicates an outlier whereby deviations by individual variables were not compensated for by offsetting movements of other variables.
${ }^{7}$ For the definition of the Borrowing ratio, see
Datastream's Company Accounts Definitions Manual,119).

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### 5.3. Dynamic ebisodes around each outlier

The search for simple dynamics could be examined by either a time series or a cross-sectional approach to our data.

The time series method focuses on the $d^{2}$ behaviour of individual firms. Episodes consisting of gradual approaches to an outlier constitute the necessarily input for failure prediction analyses. In other words, "sudden" outliers, those without any advance warning, are impossible to predict. Figure 2 on 14 shows a potentially predictable movement by firm 7 (PSUN) had observations before 1984 been included. For example, an episode consisting of years 1982 to 1984 could display a significantly different mean from that of a period running from, say, 1970 to 1981. Assuming increasing $d^{2}$ values over the 82-84 period, an even larger $d^{2}$ value (possibly an outlier) could reasonably have been anticipated.

The cross-sectional strategy centers on the distinction between years, not firms. The nonuniform distribution of the outliers over the sample period (all outliers dated from 1985 or 1984) implies a nonstationary mean. Table 7 shows the result of searching around each year for a continuous episode in which all the points were close together in Mahalanobis space and far away from the mean observation (at $1 \%$ significance).

Table 7: Search for episodes deviating from the mean state hotrllings test

| years | value | Approx.f. | sic. F. |
| :---: | :---: | :---: | :---: |
| 84 | 544.3 | 3865 | - * |
| ${ }_{84}^{84} 858$ | 5.7 | ${ }^{40.7}$ | ○* |
| 84858687 | 0.32 | 2.26 | 0.03 |
|  | 0.16 | 1.13 | 0.347 |

The decreasing $T^{2}$ statistic value suggests a strong but extremely short lived strain and long recovery period. For some reason, the mean of the observations for 1984 and for the period 84-85 appears to deviate significantly from that of the other years combined, marking a break in the pooled cross-sectional series pattern. These results reinforce our suspicion that something happened in the economy and/or in the industry during 84-85. The presence of such a break may reflect a sustained structural change in the economy. Hypothesis 4 deals with this problem.

### 5.4. Identification of Structural Changes

Figure 4 displays the outlier distribution over the sample period. All outliers occur during 1984 and 85. At first sight, the $d^{2}$ values appear to be autocorrelated over the
remaining sample period and may display significant dynamic behaviour. Given the extremely short observation period (8689), we did not test for autocorrelation.

Figure 4: Outlier Distribution over Sample Period


Though formal procedures could be applied to test for the presence of a significant break (Chow test), we believe a qualitative analysis of the printing \& publishing industry would reveal a trend reversal in 84-85. Figure 5 seems to support our feeling. The two indices stood more or less at their bottom value in 1984 and the industry's outlook improved from 1985.

Figure 5: Competitive position of Printing \& Publishing Industry


We attempted to "explain" company failure by suggesting an alternative model to the classic MDA technique. Our model, the decomposition of the Mahalanobis Distance, identifies and analyzes the actual distortions rather than its origins. Indeed, it would be futile to try to isolate causal factors because a dynamic economic system behaves very much like a nonrecursive model in which the variables cannot be arranged in any hierarchical order. Consequently, no single variable can be picked as the sole or determining cause of the state of the system afterwards.

The Mahalanobis technique incorporates some advantages over the MDA. First, there is no need to worry about the conceptualisation of the categories the observations fall into and to resort to ad hoc criteria. Second, the Mahalanobis model stands for a unique design for each year and each firm whereas the discriminant function represents an aggregate system and therefore probably incorrect. Lastly, contrary to MDA analyses, we see no need'to restrict the variables to those endogeneous, that is accounting or financial, components. External factors, such as industry and economy effects, need incorporated into the analysis.

The Mahalanobis model is, of course, not without its defects. We assumed a stable variance-covariance matrix. Our empirical test showed the presence of a nonstationary mean so the manifestation of an equally nonuniform variance cannot be dismissed. Furthermore, the predictive ability of the Mahalanobis model is restricted to those outliers that form part of longer lasting episodes. Unexpected or sudden outliers cannot be forecasted, but then how good is MDA at anticipating sudden corporate fragility ?

We believe that the comparative analysis of the contribution matrices ( $F$ ) to the distance $d^{2}$ may yield valuable insights into the failure procedure and could contribute to a more systematic approach of the corporate collapse issue. Such analysis could be best performed by a pattern recognition technique such as an artificial neural system (ANS).
Meanwhile, the Mahalanobis model is argued to assist industry regulators in the selection of "policy" variables and testing of various competing economic theories.

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# ANOTHER COUNTEREXAMPLE TO THE RANK-COLOURING CONJECTURE 

## Dirk TRAPPERS

Faculty of Applied Economics UFSIA University of Antwerp, Prinsstraat 13, B-2000 Antwerpen, Belgium

## ABSTRACT

In this paper we extend Alan and Seymour's procedure to find a counterexample of the Rank Colouring Conjecture, which was proposed by Van Nuffelen. This conjecture states that the rank of the incidence matrix of a graph is an upper bound for the chromatic number of that graph.

We assume the reader is familiar with the paper of Alon and Seymour ([1]), in which they constructed a graph with rank 29 and chromatic number 32 as a counterexample to the Rank Colouring Conjecture. By generalizing this method one can construct at least two more graphs $G$ with $\operatorname{rk}(G)<\chi(G)$, one of these is however isomorphic to the counterexample of Alon and Seymour.
The conjecture was originally put forward by Van Nuffelen in [2]. He proved the conjecture for several classes of graphs but was not able to provide a general proof for arbitrary graphs. The hope for such a proof was finally abandoned after the publication of [1]. When one looks at the results of Van Nuffelen, it remains however very plausible that the counterexamples to the conjecture should belong to a limited class of graphs. Consequently, the main motivation for this paper is to provide a narrower description of the class of counterexamples.

Alon and Seymour found their counterexample by considering Cayley graphs, these graphs are constructed by taking a finite abelian group $G$ as its vertex set and by connecting two vertices when their sum belongs to a certain predescribed subset $K$ of $G$.
When we specialize this to $G=(\mathbf{Z} / 2 \mathbf{Z})^{k}$ (the $k$-dimensional vector space over the finite field with 2 elements), we can identify subsets of $G$ with subsets of $\mathcal{P}(\{1, \ldots, k\})$ (the powerset of $\{1, \ldots, k\})$ by considering them as characteristic functions. The comnterexample of Alon and Seymour is found by putting $k=6$ and by taking $K$ as the complement in $\mathcal{P}(\{1, \ldots, k\})$ of the set

$$
\{\emptyset,\{1\},\{2\},\{3\},\{4\},\{5\},\{6\},\{1,2,3,4,5,6\}\} .
$$

We will denote $I$ for the set $\{1, \ldots, k\}$ and $\mathcal{P}(I)$ for the powerset of $I$.

Definition : We will call a set $A \subset \mathcal{P}(I)$ symmetric in the numbers $\alpha_{1}, \ldots, \alpha_{n} \in I$ if
$-\emptyset \in A$,

- every subset of $\mathcal{P}(I)$ with cardinality equal to one of the $\alpha_{i}$ is in $A$.

The set $E(H)$ in the counterexample of Alon and Seymour has $k=6, n=2$ and $\alpha_{1}=1, \alpha_{2}=6$.

We will denote the symmetric difference of two sets $X$ and $Y$ by $X * Y$, i.e.

$$
X * Y=(X \cup Y) \backslash(X \cap Y) .
$$

A subset $X$ of $I$ uniquely corresponds to a vector $\varphi_{X} \in(\mathbf{Z} / 2 \mathbf{Z})^{k}$, where $i$-th component of $\varphi_{X}$ equals 1 if $i \in X$ and equals 0 otherwise.
It is now easily seen that:

$$
\varphi_{X * Y}=\varphi_{X}+\varphi_{Y}
$$

with addition in the vector space $(\mathbf{Z} / 2 \mathbf{Z})^{k}$.
Definition : A symmetric subset $A$ of $\mathcal{P}(I)$ is called colourful if there are no $X$, $Y, Z \in A$ with $X * Y * Z=\emptyset$ or equivalently $X * Y=Z$.

Proposition 1: A colourful set $A$ cannot be symmetric in an even number $n \leq \frac{2}{3} k$.

## Proof:

Take $A$ symmetric in $n$.
Put $X=\{1, \ldots, n\}$ and $Y=\{n / 2+1, \ldots, 3 n / 2\}$, by the assumption on $n$ it follows that $X, Y \in A$.
Furthermore

$$
X * Y=\{1, \ldots, n / 2, n+1, \ldots, 3 n / 2\} \in A
$$

as the cardinality of $X * Y$ equals $n$. Consequently, $A$ is not colourful.

Proposition 2: Take $x, y, z \in I$ with $x+y=z$ and take $A$ colourful, we then have

- if $A$ is symmetric in $x$ and $y$, then $A$ is not symmetric in $z$.
- if $A$ is symmetric in $x$ and $z$, then $A$ is not symmetric in $y$,

Proof:
The proof is analogous to proposition 1.
For the first statement, take $X=\{1, \ldots, x\}$ and $Y=\{x+1, \ldots, z\}$, it follows that $X * Y=\{1, \ldots, z\}$.

And for the second statement, put $X=\{1, \ldots, x\}$ and $Z=\{1, \ldots, z\}$, then as $z>x X * Z=\{x+1, \ldots, z\}$.

Take now the hypergraph $H$ with $V(H)=I$ and $E(H)=A \subset \mathcal{P}(I)$, with $A$ symmetric.
Define the graph $G$ with vertex set equal to the vector space $V(G)=(\mathbf{Z} / 2 \mathbf{Z})^{k}$ and where $u$ and $v \in V(G)$ are adjacent if $u * v \notin E(H)$.
By demanding that $E(H)$ is colourful, we get that every stable set of $G$ has cardinality less or equal to 2 and so $\chi(G) \geq \frac{1}{2}|V(G)|=2^{k-1}$. As mentioned in [1], if $E(H) \neq\{\emptyset\}$ then in fact $\chi(G)=2^{k-1}$, as its complement contains a perfect matching.

To compute the rank of $G$, we have to count the number $t$ of $X \subset V(H)=I$ such that $|X \cap E|$ is odd for precisely $\frac{1}{2}|E(H)|$ members $E$ of $E(H)$.
By claim 2 in the paper of Alon and Seymour, it then follows that $\mathrm{rk}(G)=2^{k}-t$. To get a counterexample to the Rank-Colouring conjecture we will have to find graphs for which $t>2^{k-1}$.

To this purpose, define a function $F: I \cup\{0\} \times I \cup\{0\} \rightarrow \mathbf{Z}$ by

$$
F(a, b)=\sum_{i}\binom{b}{2 i+1}\binom{k-b}{a-2 i-1}
$$

where the summation runs over all $i$ such that $\max \{a+b-k, 1\} \leq 2 i+1 \leq$ $\min \{a, b\}$.
It is clear, by construction of $F$, that $F(a, b)$ equals the number of subsets of $I$ with cardinality $a$ that have an odd intersection with $\{1, \ldots, b\}$.

Assume further that $E(H)$ is symmetric in $\alpha_{1}, \ldots, \alpha_{n}$ then

$$
|E(H)|=1+\sum_{i=1}^{n}\binom{k}{\alpha_{i}},
$$

and so we have to find all $b \in I \cup\{0\}$ that satisfy

$$
\begin{equation*}
\sum_{i=1}^{n} F\left(\alpha_{i}, b\right)=\frac{1}{2}+\frac{1}{2} \sum_{i=1}^{n}\binom{k}{\alpha_{i}} \tag{*}
\end{equation*}
$$

If the $b_{1}, \ldots, b_{l}$ are the solutions of $(*)$, then

$$
\begin{equation*}
t=2^{k}-\operatorname{rk}(G)=\sum_{j=1}^{l}\binom{k}{b_{j}} \tag{**}
\end{equation*}
$$

It is now an easy exercise in combinatorics to verify:
Proposition 3: For $a, b \in I \cup\{0\}$ we have:

- for a even: $F(a, b)=F(a, k-b)$,
- for $a$ odd : $F(a, b)+F(a, k-b)=\binom{k}{a}$.

We then have:
Proposition 4: In order for (*) to be satisfied in a way such that $\chi(G)>\operatorname{rk}(G)$, at least one of the $\alpha_{i}$ should be even.

## Proof:

Assume all $\alpha_{i}$ odd and take $b$ a solution of (*) such that $k-b$ also satisfies (*). Now (*) and proposition 3 imply:

$$
\sum_{i=1}^{n}\binom{k}{\alpha_{i}}-\sum_{i=1}^{n} F\left(\alpha_{i}, b\right)=\frac{1}{2}+\frac{1}{2} \sum_{i=1}^{n}\binom{k}{\alpha_{i}}
$$

or

$$
\sum_{i=1}^{n} F\left(\alpha_{i}, b\right)=-\frac{1}{2}+\frac{1}{2} \sum_{i=1}^{n}\binom{k}{\alpha_{i}}
$$

Which clearly contradicts (*).
This allows us to conclude that the solution set $\left\{b_{1}, \ldots, b_{l}\right\}$ of $(*)$, can be replaced by a subset $\left\{c_{1}, \ldots, c_{l}\right\}$ of $\left\{0, \ldots,\left\lfloor\frac{k}{2}\right\rfloor\right\}$ (in (**) we may replace $b$ by $k-b$ ).
But

$$
\sum_{j=0}^{l}\binom{k}{c_{j}} \leq \sum_{i=0}^{\left\lfloor\frac{k}{2}\right\rfloor}\binom{k}{i} \leq 2^{k-1}
$$

and consequently $\mathrm{rk}(G) \geq 2^{k-1}=\chi(G)$.

Based upon these four proposition, it is possible to write a computer program, that generates sets $E(H)$ symmetric in $\alpha_{1}, \ldots, \alpha_{n}$ that satisfy:

- at least one of the $\alpha_{i}$ is even,
- all the even $\alpha_{i}$ satisfy, $\alpha_{i}>3 k / 2$,
- there are no three $\alpha_{i}, \alpha_{j}, \alpha_{k}$, such that $\alpha_{i}+\alpha_{j}=\alpha_{k}$.

The program itself can be found in Appendix 1. Appendix 2 contains the output of the program for $k=4, \ldots, 14$. On the upper part of the page, you find $F(a, b)$. Then each of the possible sets $\left\{0, \alpha_{1}, \ldots, \alpha_{n}\right\}$ are listed together with the number $g=1 / 2+1 / 2 \sum\binom{k}{\alpha_{i}}$ and $t=\sum\binom{k}{b_{j}}$.
I also verified the cases $k=15, \ldots, 18$, but these are not listed in the appendix.

The three results found are:

- $k=6, n=2$ and $\alpha_{1}=1, \alpha_{2}=6$.
- $k=6, n=2$ and $\alpha_{1}=5, \alpha_{2}=6$.
- $k=7, n=1$ and $\alpha_{1}=6$.

The corresponding sets $E(H)$ are then given by:

$$
\begin{gathered}
E\left(H_{1}\right)=\{\emptyset,\{1\},\{2\},\{3\},\{4\},\{5\},\{6\},\{1,2,3,4,5,6\}\} \\
E\left(H_{2}\right)=\{\emptyset,\{1,2,3,4,5\},\{1,2,3,4,6\},\{1,2,3,5,6\},\{1,2,4,5,6\}, \\
\{1,3,4,5,6\},\{2,3,4,5,6\},\{1,2,3,4,5,6\}\} \\
E\left(H_{3}\right)=\{\emptyset,\{1,2,3,4,5,6\},\{1,2,3,4,5,7\},\{1,2,3,4,6,7\},\{1,2,3,5,6,7\}, \\
\{1,2,4,5,6,7\},\{1,3,4,5,6,7\},\{2,3,4,5,6,7\}\}
\end{gathered}
$$

It doesn't seem very likely anymore counterexamples will be found using this method, but presently I am not able to prove this.
The first graph $G_{1}$ is the one found by Alon and Seymour, it has 64 vertices. The second one $G_{2}$ also has 64 vertices and take the map

$$
\Psi: G_{1} \rightarrow G_{2}
$$

defined by

$$
\begin{aligned}
& \Psi\left(x_{1}, x_{2}, x_{3}, x_{1}, x_{5}, x_{6}\right) \\
& \quad= \begin{cases}\left(x_{1}, x_{2}, x_{3}, x_{4}, x_{5}, x_{6}\right) & \text { if } \sum x_{i} \text { is even } \\
\left(1-x_{1}, 1-x_{2}, 1-x_{3}, 1-x_{4}, 1-x_{5}, 1-x_{6}\right) & \text { otherwise }\end{cases}
\end{aligned}
$$

It is then easily seen that $\Psi$ is an isomorphism of graphs.
The third graph $G_{3}$ is essentially different as it contains 128 vertices. We have $\operatorname{rk}\left(G_{3}\right)=128-70=58$ and $\chi\left(G_{3}\right)=64$. Altough these numbers equal the ones found for the graph we get by taking $G_{1}$ twice and joining every vertex to every vertex of the other copy (see [1]), they are not isomorphic as the complement of $G_{3}$ is a connected graph.
This proves:

Theorem : The method used by Alon and Seymour to construct a graph that doesn't satisfy the rank-colouring conjecture, only gives one other counterexample with less than 500000 vertices.

## References :

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```
# include <stdio.h>
\begin{tabular}{lll} 
\# define & K & 1 L \\
\# define & MAX & \((\) long \()(\mathrm{K}+1 \mathrm{~L})\) \\
\# define & POWER & \((\) long \()(1 \mathrm{~L} \ll(\mathrm{~K}-1 \mathrm{~L}))\) \\
\# define & TTK & \((\) long \()(2 \mathrm{~L} * \mathrm{~K} / 3 \mathrm{~L})\)
\end{tabular}
long fac[13] =
\begin{tabular}{lll}
1 L, & 1 L, & 2 L, \\
6 L, & 24 L, & 120 L, \\
720 L, & 5040 L, & 40320 L, \\
362880 L, & 3628800 L, & 39916800 L,
\end{tabular}
long \(\quad \mathrm{F}[\mathrm{MAX}][\mathrm{MAX}] ;\)
\# define \(\quad \max (\mathrm{a}, \mathrm{b}) \quad((\mathrm{a})<(\mathrm{b}) ?(\mathrm{~b}):(\mathrm{a}))\)
```

```
long combination(x,y)
```

long combination(x,y)
long x,y;
long x,y;
register long p, i
long q;
if(x < 13L)
return (fac[x] / fac[y]) / fac[x - y];
q}=\operatorname{max}(\textrm{y},\textrm{x}-\textrm{y})
p = 1L;
for(i=x;i>q;i--)
p *= i;
for(i = 1;i<= x - q;i++)
p /= i;
}
filleff()
for(a = 0;a < MAX;a++)
for(b = 0;b < MAX;b++)
sum = 0
for(j = 1;j < MAX; ; += 2)
if(j<= b \&\& j <= a \&\& a - j <= K - b)
sum += combination(b,j)
* combination( }\textrm{K}-\textrm{b},\textrm{a}-\textrm{j})\mathrm{ ;
}
}

```
long
\(\mathrm{A}[\mathrm{MAX}], \mathrm{B}[\mathrm{MAX}] ;\)
\(\operatorname{main}()\)
\{
register
long
char
FILE
sprintf(name,"graphs.\%1d",K);
if \(((\mathrm{fp}=\) fopen(name,"『")) \(==\) NULL \()\) exit(0);
filleff();
fprintf(fp,"k \(=\% 2 d \backslash n-\cdots---\backslash n \backslash n ", K) ;\)
for \((i=0 ; i<M A X \cdot i++\)
\({ }_{\text {for }}(\mathrm{j}\)
for \((\mathrm{j}=0 ; \mathrm{j}<\mathrm{MAX} ; \mathrm{j}++\) )
fprintf(fp,"\%41d ",F[i][j]);
fprintf(fp,"\n");
fprintf(fp," \(\left.{ }^{\prime} \backslash n "\right)\)
for \((\mathrm{i}=0 ; \mathrm{i}<\mathrm{MAX} ; i++\) )
\(\mathrm{A}[\mathrm{i}]=0\);
next:
\(\begin{aligned} \text { for }(\mathrm{i}= & 1 ; \mathrm{i}<=\text { MAX } ; i++) \\ & \{ \\ & \text { if }(\mathrm{i}==\text { MAX })\end{aligned}\)
\{ \(\mathrm{flose}(\mathrm{fp})\); exit(0);
\(\begin{aligned} \operatorname{if}(\mathrm{A}[\mathrm{i}]= & \left.\frac{\}}{=} 0\right) \\ & \{ \\ & \mathrm{A}[\mathrm{i}]=1 ;\end{aligned}\)
break;
\({ }_{0}{ }^{\text {brea }}\)
\(\mathrm{A}[\mathrm{i}]=0\);
for \((\mathrm{i}=\stackrel{\}}{1 ; i<M A X ; i++)} \underset{i f(A[i] \& \& i \%}{ }\)
if(A[i] \&\& i \% 2 \(==0\) )
if( \(\mathrm{i}<=\) TTK)
goto next;
goto even
/* A should contain at least 1 even number, but if \(2^{*} x<=T T K\) then \(2^{*} x\) is excluded */
goto next;
even:
for \((\mathrm{i}=1 ; \mathrm{i}<\) MAX; \(\mathrm{i}++\) ) for \(\left(\mathrm{j}=1 ; \mathrm{j}<\mathrm{MAX}_{\mathrm{j}} \mathrm{j}++\right.\) )
main
\[
\begin{aligned}
& \{ \\
& \text { if(i }+\mathrm{j}>=\text { MAX }) \\
& \text { continue; } \\
& \text { if( } \mathrm{A}[\mathrm{i}] \& \& A[\mathrm{j}] \& \& \mathrm{~A}[\mathrm{i}+\mathrm{j}]) \\
& \text { goto next; } \\
& \} \quad
\end{aligned}
\]
```

number $=0$;
for $(\mathrm{i}=1 ; \mathrm{i}<$ MAX; $;++$ )
if(A[i])

```
if(number \(\% 2=0\) )
        goto next;
number++;
number \(/=2\);
for \((\mathrm{i}=0 ; \mathrm{i}<\mathrm{MAX} ; \mathrm{i}++\) )
        \(\mathrm{B}[\mathrm{i}]=0\);
for( \(\mathrm{i}=1 ; \mathrm{i}<\mathrm{MAX} ; i++\) )
        if(A[i])
                for \(\left(j=0 ; j<\operatorname{MAX}_{j} \mathrm{j}++\right.\) )
                    \(\mathrm{B}[\mathrm{j}]+=\mathrm{F}[\mathrm{i}][\mathrm{j}] ;\)
sum \(=0\);

\(\begin{aligned} \text { for }(\mathrm{i}= & 0 ; \mathrm{i}<\mathrm{MAX} ; \mathrm{i}++) \\ & \mathbf{i f}(\mathrm{B}[\mathrm{i}]==\text { number })\end{aligned}\)
                sum \(+=\) combination \((K, i)\) :
fprintf(fp,"\%2d ",0);
for \((\mathrm{i}=1 ; \mathrm{i}<\mathrm{MAX} ; \mathrm{i}++\) )
        if(A[i])
            fprintf(fp,"\%21d ",i);
        else
            fprintf(fp," ");
fprintf(fp," \(\rightarrow\) ( \(\mathrm{g}=\% 61 \mathrm{~d}\) ) \%61d",number.sum);
if(sum > POWER)
        fprintf(fp," !!!");
fprintf(fp,"\n");
goto next;
\[
\begin{aligned}
& \begin{array}{l}
k=4 \\
-
\end{array} \\
& \begin{array}{lllll}
0 & 0 & 0 & 0 & 0 \\
0 & 1 & 2 & 3 & 4 \\
0 & 3 & 4 & 3 & 0 \\
0 & 3 & 2 & 1 & 4 \\
0 & 1 & 0 & 1 & 0
\end{array} \\
& \begin{array}{lllllll}
0 & & & 4 & -\infty & (g= & 1) \\
0 & 1 & 3 & --> & (g= & 3) & 0 \\
0 & 3 & --> & (g= & 3) & 0 \\
k=5 & & & & & &
\end{array} \\
& \begin{array}{rrrrrr}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 2 & 3 & 4 & 5 \\
0 & 4 & 6 & 6 & 4 & 0 \\
0 & 6 & 6 & 4 & 4 & 10 \\
0 & 4 & 2 & 2 & 4 & 0 \\
0 & 1 & 0 & 1 & 0 & 1
\end{array} \\
& \begin{array}{llllllr}
0 & & 4 & -> & (g= & 3) & 0 \\
0 & 3 & 4 & -> & (g= & 8) & 15
\end{array} \\
& k=6 \\
& \begin{array}{rrrrrrr}
0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 2 & 3 & 4 & 5 & 6 \\
0 & 5 & 8 & 9 & 8 & 5 & 0 \\
0 & 10 & 12 & 10 & 8 & 10 & 20 \\
0 & 10 & 8 & 6 & 8 & 10 & 0 \\
0 & 5 & 2 & 3 & 4 & 1 & 6 \\
0 & 1 & 0 & 1 & 0 & 1 & 0
\end{array} \\
& \begin{array}{llllllll}
0 & & 6 & --> & (9= & 1) & 32 \\
0 & 1 & & 6 & --> & (9= & 4) & 35!!! \\
0 & 5 & 6 & --> & (9= & 4) & 35!!! \\
k=7 & & & & &
\end{array} \\
& \begin{array}{rrrrrrrr}
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 2 & 3 & 4 & 5 & 6 & 7 \\
0 & 6 & 10 & 12 & 12 & 10 & 6 & 0 \\
0 & 15 & 20 & 19 & 16 & 15 & 20 & 35 \\
0 & 20 & 20 & 16 & 16 & 20 & 20 & 0 \\
0 & 15 & 10 & 9 & 12 & 11 & 6 & 21 \\
0 & 6 & 2 & 4 & 4 & 2 & 6 & 0 \\
0 & 1 & 0 & 1 & 0 & 1 & 0 & 1
\end{array} \\
& \begin{array}{lrllllll}
0 & & 6 & & & (g= & 4) & 70 \\
0 & 5 & 6 & 7 & --> & (g= & 15) & 0
\end{array} \\
& \begin{array}{rrrrrrrrr}
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 \\
0 & 7 & 12 & 15 & 16 & 15 & 12 & 7 & 0 \\
0 & 21 & 30 & 31 & 28 & 25 & 26 & 35 & 56 \\
0 & 35 & 40 & 35 & 32 & 35 & 40 & 35 & 0 \\
0 & 35 & 30 & 25 & 28 & 31 & 26 & 21 & 56 \\
0 & 21 & 12 & 13 & 16 & 13 & 12 & 21 & 0 \\
0 & 7 & 2 & 5 & 4 & 3 & 6 & 1 & 8 \\
0 & 1 & 0 & 1 & 0 & 1 & 0 & 1 & 0
\end{array} \\
& 0 \\
& 8-->\text { ( } g \\
& \text { 1) } 128
\end{aligned}
\]

\(k=10\)


\begin{tabular}{|c|c|c|c|c|c|c|c|c|c|c|}
\hline 0 & & & & & 8 & & --> (g & \(=\) & 83) & 0 \\
\hline 0 & 1 & 3 & & & 8 & & --> (g & & 171) & 0 \\
\hline 0 & & & 5 & & 8 & & --> 1g & & 314) & 0 \\
\hline 0 & & & & 7 & 8 & & --> (g & & 248) & 0 \\
\hline 0 & & & 5 & 7 & 8 & & --> 19 & & 479) & 0 \\
\hline 0 & & 3 & & & 8 & 9 & --> (g & = & 193) & 0 \\
\hline 0 & & 3 & & 7 & 8 & 9 & --> 19 & & 358) & 0 \\
\hline 0 & & & & & & 10 & --> (g & = & 6) & 924 \\
\hline 0 & 1 & 3 & & & & 10 & --> \(\mathrm{g}^{\text {g }}\) & & 94) & 0 \\
\hline 0 & & & & 7 & & 10 & --> (g & & 171) & 0 \\
\hline 0 & 1 & & & & 8 & 10 & --> (g & & 94) & 0 \\
\hline
\end{tabular}

\(k=12\)



\(\mathrm{k}=13\)
\begin{tabular}{rrrrrrrrrrrrrr}
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 & 11 & 12 & 13 \\
0 & 12 & 22 & 30 & 36 & 40 & 42 & 42 & 40 & 36 & 30 & 22 & 12 & 0 \\
0 & 66 & 110 & 136 & 148 & 150 & 146 & 140 & 136 & 138 & 150 & 176 & 220 & 286 \\
0 & 220 & 330 & 370 & 372 & 360 & 350 & 350 & 360 & 372 & 370 & 330 & 220 & 0 \\
0 & 495 & 660 & 675 & 648 & 631 & 636 & 651 & 656 & 639 & 612 & 627 & 792 & 1287 \\
0 & 792 & 924 & 876 & 840 & 848 & 868 & 868 & 848 & 840 & 876 & 924 & 792 & 0 \\
0 & 924 & 924 & 840 & 840 & 868 & 868 & 848 & 848 & 876 & 876 & 792 & 792 & 1716 \\
0 & 792 & 660 & 612 & 648 & 656 & 636 & 636 & 656 & 648 & 612 & 660 & 792 & 0 \\
0 & 495 & 330 & 345 & 372 & 355 & 350 & 365 & 360 & 343 & 370 & 385 & 220 & 715 \\
0 & 220 & 110 & 150 & 148 & 136 & 146 & 146 & 136 & 148 & 150 & 110 & 220 & 0 \\
0 & 66 & 22 & 48 & 36 & 38 & 42 & 36 & 40 & 42 & 30 & 56 & 12 & 78 \\
0 & 12 & 2 & 10 & 4 & 8 & 6 & 6 & 8 & 4 & 10 & 2 & 12 & 0 \\
0 & 1 & 0 & 1 & 0 & 1 & 0 & 1 & 0 & 1 & 0 & 1 & 0 & 1
\end{tabular}
\begin{tabular}{|c|c|c|c|c|c|c|c|c|c|c|c|c|c|}
\hline 0 & 1 & & & & & 10 & & & & --> (g & \(=\) & 150) & 0 \\
\hline 0 & 1 & 3 & & & & 10 & & & & --> (g & & 293) & 1716 \\
\hline 0 & 1 & & & 7 & & 10 & & & & --> (g & & 1008) & 0 \\
\hline 0 & & & & & 9 & 10 & & & & --> (g & & 501) & 0 \\
\hline 0 & & 3 & & & 9 & 10 & & & & --> (g & & 644) & 0 \\
\hline 0 & & & & 7 & 9 & 10 & & & & --> (g & & 1359) & 3003 \\
\hline 0 & & & & & 9 & 10 & 11 & & & --> (g & & 540) & 0 \\
\hline 0 & & 3 & & & 9 & 10 & 11 & & & \(\cdots 19\) & & 683) & 0 \\
\hline 0 & & & & 7 & 9 & 10 & 11 & & & --> (g & & 1398) & 0 \\
\hline 0 & & & & & & & & 12 & & --> (g & & 7) & 0 \\
\hline 0 & & 3 & & & & & & 12 & & --> (g & & 150) & 0 \\
\hline 0 & 1 & & 5 & & & & & 12 & & --> (g & & 657) & 0 \\
\hline 0 & 1 & 3 & 5 & & & & & 12 & & --> (g & & 800) & 0 \\
\hline 0 & & & & & & & & 12 & & --> ( \(g\) & & 865) & 0 \\
\hline 0 & & 3 & & 7 & & & & 12 & & --> (g & & 1008) & 0 \\
\hline 0 & 1 & & & & 9 & & & 12 & & --> (g & & 371) & 0 \\
\hline 0 & & & 5 & & 9 & & & 12 & & --> 19 & & 1008) & 0 \\
\hline 0 & 1 & & & 7 & 9 & & & 12 & & --> (9) & & 1229) & 0 \\
\hline 0 & & & & & & 10 & & 12 & & \(\cdots 19\) & & 150) & 0 \\
\hline 0 & & 3 & & & & 10 & & 12 & & --> (g & & 293) & 0 \\
\hline 0 & & & & 7 & & 10 & & 12 & & --> (g & & 1008) & 0 \\
\hline 0 & & & & & & & 11 & 12 & & --> 19 & & \(46)\) & 2002 \\
\hline 0 & & 3 & & & & & & & & --> (g & & 189) & 0 \\
\hline 0 & & & & 7 & & & 11 & & & --> (g & \(=\) & 904) & 0 \\
\hline 0 & & 3 & & 7 & & & & & & \(-->\) ( g & & 1047) & 0 \\
\hline 0 & & & 5 & & 9 & & 11 & & & \(\cdots\) (g & & 1047) & 0 \\
\hline 0 & & & & & & 10 & 11 & 12 & & --> (g & & 189) & 0 \\
\hline 0 & & 3 & & & & 10 & 11 & 12 & & --> (g & \(=\) & 332) & 2002 \\
\hline 0 & & & & 7 & & 10 & 11 & 12 & & --> (g & & 1047) & 0 \\
\hline 0 & & & & & & 10 & & & 13 & --> (g & \(=\) & 144) & 0 \\
\hline 0 & & & & 7 & & 10 & & & 13 & --> (g & & 1002) & 0 \\
\hline 0 & & & & & & 10 & 11 & & 13 & --> 1 g & \(=\) & 183) & 1716 \\
\hline 0 & & & & 7 & & 10 & 11 & & 13 & --> 19 & & 1041) & 0 \\
\hline 0 & & & 5 & & & & & 12 & 13 & --> 19 & & 651) & 0 \\
\hline 0 & & 3 & 5 & & & & & 12 & 13 & \(\rightarrow\) ( 9 & & 794) & 0 \\
\hline 0 & & & & & 9 & & & 12 & 13 & --> (g & & 365) & 0 \\
\hline 0 & & & & 7 & 9 & & & 12 & 13 & --> (g & & 1223) & 0 \\
\hline 0 & & & & & 9 & 10 & & 12 & 13 & \(\cdots\) (g & & 508) & 0 \\
\hline 0 & & & & 7 & 9 & 10 & & 12 & 13 & --> (g & & 1366) & 1794 \\
\hline 0 & & & 5 & & & & 11 & 12 & 13 & --> 19 & & 690) & 0 \\
\hline 0 & & 3 & 5 & & & & 11 & 12 & 13 & --> 19 & & 833) & 0 \\
\hline 0 & & & & & 9 & & 11 & 12 & 13 & --> 19 & & 404) & 286 \\
\hline
\end{tabular}


0000000000000000000000000000000000000000000000000000000000000000 1

\begin{tabular}{|c|c|}
\hline &  \\
\hline  &  \\
\hline  & nNenNen nNmN \\
\hline
\end{tabular}





This is a picture of the complementary graph of a graph \(G\) that looks very much like the complementary graph of \(G_{1}\) (and of \(G_{2}\) ). To get a real picture of the complement of \(G_{1}\), you should also connect the vertices that are diametrically placed with respect to the central point of the drawing.
In a drawing of the complement of \(G_{1}\), one should label the vertices as follows:
\begin{tabular}{llllllll}
00 & 01 & 02 & 03 & 04 & 05 & 06 & 07 \\
10 & 11 & 12 & 13 & 14 & 15 & 16 & 17 \\
20 & 21 & 22 & 23 & 24 & 25 & 26 & 27 \\
30 & 31 & 32 & 33 & 34 & 35 & 36 & 37 \\
40 & 41 & 42 & 43 & 44 & 45 & 46 & 47 \\
50 & 51 & 52 & 53 & 54 & 55 & 56 & 57 \\
60 & 61 & 62 & 63 & 64 & 65 & 66 & 67 \\
70 & 71 & 72 & 73 & 74 & 75 & 76 & 77
\end{tabular}

The labels are the octal numbers that correspond to the binary representation in \((\mathbb{Z} / 2 \mathbb{Z})^{6}\).
E.g.: The octal number 53 corresponds with the binary number 101011 and with the vector \((1,0,1,0,1,1)\) of \((\mathbb{Z} / 2 \mathbb{Z})^{6}\).

To get a picture of the complement of \(G_{2}\), one should start from the same graph but choose the labels:
\begin{tabular}{llllllll}
00 & 76 & 75 & 03 & 73 & 05 & 06 & 70 \\
67 & 11 & 12 & 64 & 14 & 62 & 61 & 17 \\
57 & 21 & 22 & 54 & 24 & 52 & 51 & 27 \\
30 & 46 & 45 & 33 & 43 & 35 & 36 & 40 \\
37 & 41 & 42 & 34 & 44 & 32 & 31 & 47 \\
50 & 26 & 25 & 53 & 23 & 55 & 56 & 20 \\
60 & 16 & 15 & 63 & 13 & 65 & 66 & 10 \\
07 & 71 & 72 & 04 & 74 & 02 & 01 & 77
\end{tabular}


\section*{Review of the book: "Histoire de la Statistique"} by J.J. DROESBEKE \& Ph TASSI, P.U.F., Paris,

Collection "Que Sais-je ?" n \({ }^{\circ}\) 2527, 1990

It is a challenge to write the history of a scientific discipline in the 128 pages of a "Que Sais-je ?". Though, that is what J.J. DROESBEKE and Ph. TASSI have done for Statistics. This little book is pleasant to read and it will give the reader (with some previous knowledge in the field) a general view of the roots and main streams of development of Statistical Theory. This book is all the more welcome that French litterature is not so rich in histories of Statistics.

The first chapter deals with descriptive statistics. It is clear that Statistics is born of the need and desire of summarizing and interpreting large sets of collected data. Hence, graphical representations (first known bar chart in 1786), typical values (means, median, variance since Gauss and the mean squares method, ...), curve fitting, correlation and regression. Special attention is paid to the index numbers born three centuries ago (in England) for synthesizing price evolution.

Although many descriptive indices can be defined and used without probabilistic assumptions, probability is present in Statistics roughly since the mid XVIIIth century with the theory of errors of observation. As the main pieces of statistical theory rely heavily on Probability, the authors give a brief account of its history in the second chapter : a 16 pages trip trom Pascal, Fermat and the Bernouilli Brothers to Kolmogorov. One learns for instance that the normal law, often attributed to Laplace and Gauss, can be traced back to the Ars Conjectandi of Jacques Bernotiili (1654-1705).

The main body of statistical knowledge is then split under the five following headings: Survey Sampling (Les sondages), The rise of statistical inference, Non parametric Statistics and robustness, Time series analysis, Data analysis. One may regret (but there is always something to regret in such a work) the absence of a chapter devoted to subjects like the general linear model and the analysis of variance or to the design of experiments. Anyway.

The fact that Survey Sampling appears first in the main chapters is probably due to its link with census (recensements) which are very ancient (perhaps 5000 B.C.). From the seventieth century on, the governments have shown growing interest in collecting data about the people and country they rule and, of course, at the cheapest possible cost. This led to survey sampling. The authors mainly depict the difficult genesis of the notion of representative sample.

The chapter on estimation and tests in centered on the competing theories of Fisher on one hand and Neyman-Pearson on the other hand. My regret here is that the respective positions of the bayesian and non-bayesian statisticians are not more extensively outlined. The authors would probably argue that the debate is not sufficiently mature by now to be able to give an objective presentation.

The next chapter is concerned with two types of more recently developed techniques : the theory of robust statistics and non-parametric statistics.

The fact that the observations depend on a very special and familiar parameter called time singularizes Time series analysis among stochastic processes and justifies special attention. The methods reviewed include graphical approaches, analysis in the frequency domain (based on the periodogram), methods of decomposition (like the famous CENSUS), analysis in the time domain (based on the correlograms), stochastic models including BoxJenkins methodology.

The last but one chapter - the last one being devoted to a short biography of eight of the most prominent statisticians, those who made Statistics what it is is about data analysis in the French meaning of the word. This chapter is "a gift of B. Fichet" to the authors. The history of the numerous methods is outlined : "metric methods" as principal component analysis, correspondance analysis and discriminant analysis (with or without distributional, i.e. normal, assumptions); non-metric methods, i.e. essentially, clustering.

In conclusion, the reading of this small book is recommanded, as the authors say, to teachers of, searchers in and users of statistical methods Several reading levels are possible : the general lines will be perceived by everyone while the details, especially in recent theories, will be understood by people with higher degree of education in the field.

\footnotetext{
Marc PIRLOT,
Faculté Polytechnique de Mons.
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[^0]:    1 Donaldson's (1969) Financial Mobility Strategy attempted to integrate unanticipated events and cash flow patterns in a systematic manner.

    2 Or any other event, such as economic instability (Baestaens,1990)
    ${ }^{3}$ Multiple Discriminant Analysis

[^1]:    613 firms times 6 years ( 84 to 89 ) +4 firms times 1 Year (since these firms report in April, accounting information for 1990 was available at the time of our data collection).

