Mémoire D'HABILITATION À DIRIGER DES RECHERCHES

Mention Mathématiques et Applications

par

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TITRE DE LA THÈSE :

Contributions à l'étude des systèmes d'équations implicites dans les anneaux de séries

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1.1. Listes des travaux présentés pour l’habilitation

- [Ron13b] About the algebraic closure of the field of power series in several variables in characteristic zero, prépublication.

1.2. Publications découlant de ma thèse de doctorat


1.3. Autres travaux

• [Ro-Sp13] (Avec Mark Spivakovsky) The analogue of Izumi’s Theorem for Abhyankar valuations, prépublication.
Nous présentons dans ce mémoire une partie des résultats que nous avons obtenus depuis 2006. Ces travaux portent sur l’étude des solutions \( y = y(x) \) d’équations (ou de systèmes d’équations) du type
\[
f(x, y) = 0
\]
où \( f(x, y) \) appartient à une classe de (germes) de fonctions vérifiant le théorème des fonctions implicites, et \( x \) et \( y \) sont des multivariables. Les types de classes de fonctions considérées sont, par exemple, les germes de fonctions analytiques à l’origine de \( \mathbb{C}^{n+m} \), les séries algébriques ou les séries formelles en \( n + m \) variables, etc.
Le premier résultat général sur ce type d’équations est bien évidemment le théorème des fonctions implicites qui nous donne l’existence de solutions à ce type d’équations sous une hypothèse de lissité du système par rapport aux variables \( y \). Le cas singulier est donc celui qui nous intéresse. Le second résultat le plus connu est le théorème suivant dû à M. Artin :

**Théorème 2.0.1.** — [Ar68] Les solutions \( y = y(x) \) (germes de fonctions) analytiques d’un système d’équations analytiques de la forme \( f(x, y) = 0 \) sont denses dans l’ensemble des solutions formelles.

Ici la densité signifie que pour toute solution formelle il existe une solution analytique dont le développement de Taylor coïncide avec le développement de Taylor de la solution formelle à un ordre arbitrairement grand.

Les questions abordées dans ce mémoire sont des trois types suivants :

- Pour construire des solutions analytiques à un système \( f(x, y) = 0 \), M. Artin a montré qu’il suffit de construire des solutions approchées à un ordre grand.
Le problème est d’avoir des bornes effectives sur ces ordres d’approximation.

- Le deuxième problème consiste à savoir dans quelle mesure le théorème de Artin cité ci-dessus reste valable quand on impose des contraintes sur les solutions (contraintes du type : certains des \( y_i \) ne dépendent que de certaines des variables \( x_j \)).

- Le troisième consiste à comprendre et décrire la structure de l’ensemble des solutions d’un tel système.

Ce mémoire comprend deux parties : la première présente un survol des résultats connus sur l’approximation de Artin, thème qui a connu de grands développements dans les années 70 et 80 et pour lequel il n’existe pas de texte de référence récent et complet (on peut néanmoins citer [Ra69], [Te94], [Po00]). Néanmoins prétendre être complet serait une gageure, et nous nous restreignons à un cadre très algébrique (en particulier nous ne parlons pas d’approximation globale dans les variétés de Nash ou d’approximation en géométrie CR). C’est pourquoi nous l’avons rédigée en anglais car beaucoup de personnes susceptibles d’être intéressées par ce texte ne lisent pas le français. La seconde partie présente avec plus de détails nos contributions à ces trois problèmes :

- Pour la première question, nous étudions le cas des systèmes d’équations polynomiales \( f(x, y) = 0 \) où \( x = (x_1, x_2) \). Ceci correspond aux articles [Ron10a] et [Ron13a].

- Pour la seconde question, nous étudions le cas de systèmes d’équations "cyclindriques" ("nested" en anglais), i.e. le cas où \( y_i = y_i(x_1, ..., x_{n(i)}) \) et \( (n(i))_i \) est une suite croissante. Dans notre travail ces systèmes d’équations proviennent d’équations "fonctionnelles" du type \( f \circ g(y) = h(y) \) où \( g(y) \) et \( h(y) \) sont des germes de fonctions analytiques (par exemple) et l’inconnue \( f(x) \) est une série formelle. Ceci correspond aux articles [Ron08a], [Ron09b].

- Pour la dernière question, nous étudions le cas d’une équation polynomiale en une seule variable \( y \) à coefficients dans l’anneau des séries formelles en \( n \) variables sur un corps de caractéristique nulle. C’est-à-dire nous nous intéressons à la clôture algébrique du corps des séries formelles en plusieurs variables en caractéristique nulle. Ceci correspond aux articles [PaRo12] et [Ron13b].
3.1. Introduction

The aim of this chapter is to present the Artin Approximation Theorem and some related results. The problem we are interested in is to find analytic solutions of some systems of equations when this system admits formal power series solutions and the Artin Approximation Theorem yields a positive answer to this problem. We begin this chapter by giving several examples explaining what this sentence means exactly. Then we will present the state of the art on this problem. There are essentially three parts: the first part is devoted to present the Artin Approximation Theorem and its generalizations; the second part presents a stronger version of Artin Approximation Theorem; the last part is mainly devoted to explore the Artin Approximation Problem in the case of constraints. Sections 3.6, 3.7 and 3.8 present the algebraic material used in this chapter (Weierstrass Preparation Theorem, excellent rings, étales morphisms and Henselian rings).

We do not give the proofs of all the results presented in this chapter but, at least, we always try to outline the proofs and give the main arguments.

Example 3.1.1. — Let us consider the following curve \( C := \{(t^3, t^4, t^5), t \in \mathbb{C}\} \) in \( \mathbb{C}^3 \). This curve is an algebraic set which means that it is the zero locus of polynomials in three variables. Indeed, we can check that \( C \) is the zero locus of the polynomials \( f := y^2 - xz \), \( g := yz - x^3 \) and \( h := z^2 - x^2y \). If we consider the zero locus of any two of these polynomials we always get a set larger than \( C \). The complex dimension of the zero locus of one non-constant polynomial in three variables is 2 (such a set is called a hypersurface of \( \mathbb{C}^3 \)). Here \( C \) is the intersection of the zero locus of three hypersurfaces and not of two of them, but its complex dimension is 1.

In fact we can see this phenomenon as follows: we call an algebraic relation between \( f, g \) and \( h \) any element of the kernel of the linear map \( \varphi : \mathbb{C}[x, y, z]^3 \rightarrow \)
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\( \mathbb{C}[x, y, z] \) defined by \( \varphi(a, b, c) := af + bg + ch \). Obviously \( r_1 := (g, -f, 0) \), \( r_2 := (h, 0, -f) \) and \( r_3 := (0, h, -g) \in \text{Ker}(\varphi) \). These are called the trivial relations between \( f, g \) and \( h \). But in our case there is one more relation which is \( r_0 := (z, y, -x) \) and \( r_0 \) cannot be written as \( a_1r_1 + a_2r_2 + a_3r_3 \) with \( a_1, a_2 \) and \( a_3 \in \mathbb{C}[x, y, z] \), which means that \( r_0 \) is not in the sub-\( \mathbb{C}[x, y, z] \)-module of \( \mathbb{C}[x, y, z] \) generated by \( r_1, r_2 \) and \( r_3 \).

On the other hand we can prove that \( \text{Ker}(\varphi) \) is generated by \( r_0, r_1, r_2 \) and \( r_3 \).

Let \( X \) be the common zero locus of \( f \) and \( g \). If \( (x, y, z) \in X \) and \( x \neq 0 \), then \( h = \frac{zf + yg}{x} = 0 \) thus \( (x, y, z) \in C \). If \( (x, y, z) \in X \) and \( x = 0 \), then \( y = 0 \). Geometrically this means that \( X \) is the union of \( C \) and the \( z \)-axis, i.e. the union of two curves.

Now let us denote by \( \mathbb{C}[x, y, z] \) the ring of formal power series with coefficients in \( \mathbb{C} \). We can also consider formal relations between \( f, g \) and \( h \), that is elements of the kernel of the map \( \mathbb{C}[x, y, z]^3 \to \mathbb{C}[x, y, z] \) induced by \( \varphi \). Any element of the form \( a_0r_0 + a_1r_1 + a_2r_2 + a_3r_3 \) is a formal relation as soon as \( a_0, a_1, a_2, a_3 \in \mathbb{C}[x, y, z] \).

In fact any formal relation is of this form, i.e. the algebraic relations generate the formal and analytic relations. We can show this as follows: we can assign a formal relation then we can write \( \exists \sum \in \mathbb{N} \sum_i a_i b_i + c_i \in \mathbb{C}[x, y, z] \).

with the assumption \( b_i = c_i = 0 \) for \( i < 0 \). Thus \( (a_0, 0, 0) \), \( (a_1, b_0, 0) \) and any \( (a_i, b_{i-1}, c_{i-2}) \) for \( 2 \leq i \), are in \( \text{Ker}(\varphi) \), thus are homogeneous combinations of \( r_0, r_1, r_2 \) and \( r_3 \). Hence \( (a, b, c) \) is a combination of \( r_0, r_1, r_2 \) and \( r_3 \) with coefficients in \( \mathbb{C}[x, y, z] \).

Now we can investigate the same problem by replacing the ring of formal power series by \( \mathbb{C}\{x, y, z\} \), the ring of convergent power series with coefficients in \( \mathbb{C} \), i.e.

\[
\mathbb{C}\{x, y, z\} := \left\{ \sum_{i,j,k} a_{i,j,k} x^i y^j z^k / \exists \rho > 0, \sum_{i,j,k} |a_{i,j,k}| \rho^{i+j+k} < \infty \right\}
\]

We can also consider analytic relations between \( f, g \) and \( h \), i.e. elements of the kernel of the map \( \mathbb{C}\{x, y, z\}^3 \to \mathbb{C}\{x, y, z\} \) induced by \( \varphi \). From the formal case we see that any analytic relation \( r \) is of the form \( a_0r_0 + a_1r_1 + a_2r_2 + a_3r_3 \).
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with $a_i \in \mathbb{C}[x, y, z]$ for $0 \leq i \leq 4$. In fact we can prove that $a_i \in \mathbb{C}\{x, y, z\}$ for $0 \leq i \leq 4$. Let us remark that, saying that $r = a_0r_0 + a_1r_1 + a_2r_2 + a_3r_3$ is equivalent to say that $a_0, \ldots, a_3$ satisfy a system of three affine equations with analytic coefficients. This is the first example of the problem we are interested in: if we some equations with analytic coefficients have formal solutions do they have analytic solutions? Artin Approximation Theorem yields an answer to this problem. Here is the first theorem proven by M. Artin in 1968:

**Theorem 3.1.2 (Artin Approximation Theorem)**

Let $f(x, y)$ be a vector of convergent power series over $\mathbb{C}$ in two sets of variables $x$ and $y$. Assume given a formal power series solution $\hat{y}(x)$,

$$f(x, \hat{y}(x)) = 0.$$

Then there exists, for any $c \in \mathbb{N}$, a convergent power series solution $y(x)$,

$$f(x, y(x)) = 0$$

which coincides with $\hat{y}(x)$ up to degree $c$,

$$y(x) \equiv \hat{y}(x) \ mod (x)^c.$$

We can define a topology on $\mathbb{C}[x]$ by saying that two power series are close if their difference is in a high power of the maximal ideal $(x)$. Thus we can reformulate Theorem 3.1.2 as: formal power series solutions of a system of analytic equations may be approximated by convergent power series solutions.

**Example 3.1.3.** — A special case of Theorem 3.1.2 and a generalization of Example 3.1.1 occurs when $f$ is linear in $y$, say $f(x, y) = \sum f_i(x)y_i$, where $f_i(x)$ is a vector of convergent power series with $r$ coordinates for any $i$. A solution $y(x)$ of $f(x, y) = 0$ is a relation between the $f_i(x)$. In this case the formal relations are linear combinations of analytic combinations with coefficients in $\mathbb{C}[x]$. In term of commutative algebra, this is expressed as the flatness of the ring of formal power series over the ring of convergent powers series, a result which can be proven via the Artin-Rees Lemma.

It means that if $\hat{y}(x)$ is a formal solution of $f(x, y) = 0$, then there exist analytic solutions of $f(x, y) = 0$ denoted by $\tilde{y}_i(x)$, $1 \leq i \leq s$, and formal power series $\hat{b}_1(x), \ldots, \hat{b}_s(x)$, such that $\hat{y}(x) = \sum_i \hat{b}_i(x)\tilde{y}_i(x)$. Thus, by replacing in the previous sum the $\hat{b}_i(x)$ by their truncation at order $c$, we obtain an analytic solution of $f(x, y) = 0$ coinciding with $\hat{y}(c)$ up to degree $c$. If the $f_i(x)$’s are vectors of polynomials then the formal relations are also linear combinations of algebraic relations since the ring of formal power series is flat over the ring of polynomials, and Theorem 3.1.2 remains true if $f(x, y)$ is linear in $y$ and $\mathbb{C}\{x\}$ is replaced by $\mathbb{C}[x]$. 
Example 3.1.4. — A slight generalization of the previous example is when \( f(x, y) \) is a vector of polynomials in \( y \) of degree one with coefficients in \( \mathbb{C}[x] \) (resp. \( \mathbb{C}[x] \), say
\[
 f(x, y) = \sum_{i=1}^{m} f_i(x)y_i + b(x)
\]
where the \( f_i(x) \)'s and \( b(x) \) are vectors of convergent power series (resp. polynomials). Here \( x \) and \( y \) are multi-variables If \( \tilde{y}(x) \) is a formal power series solution of \( f(x, y) = 0 \), then \((\tilde{y}(x), 1)\) is a formal power series solution of \( g(x, y, z) = 0 \) where
\[
 g(x, y, z) := \sum_{i=1}^{m} f_i(x)y_i + b(x)z
\]
and \( z \) is a single variable. Thus using the flatness of \( \mathbb{C}[x] \) over \( \mathbb{C}[x] \) (resp. \( \mathbb{C}[x] \)) (Example 3.1.3), we can approximate \((\tilde{y}(x), 1)\) by a convergent power series (resp. polynomial) solution \((\tilde{y}(x), \tilde{z}(x))\) which coincides with \((\tilde{y}(x), 1)\) up to degree \( c \). In order to obtain a solution of \( f(x, y) = 0 \) we would like to be able to divide \( \tilde{y}(x) \) by \( \tilde{z}(x) \) since \( \tilde{y}(x)\tilde{z}(x)^{-1} \) would be a solution of \( f(x, y) = 0 \) approximating \( \tilde{y}(x) \). We can remark that, if \( c \geq 1 \), then \( \tilde{z}(0) = 1 \) thus \( \tilde{z}(x) \) is not in the ideal \((x)\). But \( \mathbb{C}(x) \) is a local ring. We call a local ring any ring \( A \) that has only one maximal ideal. This is equivalent to say that \( A \) is the disjoint union of one ideal (its only maximal ideal) and of the set of units in \( A \). In particular \( \tilde{z}(x)^{-1} \) is invertible in \( \mathbb{C}(x) \), hence we can approximate formal power series solutions of \( f(x, y) = 0 \) by convergent power series solutions.

In the case \((\tilde{y}(x), \tilde{z}(x))\) is a polynomial solution of \( g(x, y, z) = 0 \), \( \tilde{z}(x) \) is not invertible in general in \( \mathbb{C}[x] \) since it is not a local ring. For instance set
\[
 f(x, y) := (1 - x)y - 1
\]
where \( x \) and \( y \) are single variables. Then \( y(x) := \sum_{n=0}^{\infty} x^n = \frac{1}{1 - x} \) is the only formal power series solution of \( f(x, y) = 0 \), but \( y(x) \) is not a polynomial. Thus we cannot approximate the roots of \( f \) in \( \mathbb{C}[x] \) by roots of \( f \) in \( \mathbb{C}[x] \).

But instead of working in \( \mathbb{C}[x] \) we can work in \( \mathbb{C}[x]_{(x)} \) which is the ring of rational functions whose denominator does not vanish at 0. This ring is a local ring. Since \( \tilde{z}(0) \neq 0 \), then \( \tilde{y}(x)\tilde{z}(x)^{-1} \) is a vector of rational function of \( \mathbb{C}[x]_{(x)} \). In particular any system of polynomial equations of degree one with coefficients in \( \mathbb{C}[x] \) which has solutions in \( \mathbb{C}[x] \) has solutions in \( \mathbb{C}[x]_{(x)} \).

Example 3.1.5. — The next example we are looking at is the following: set \( f \in A \) where \( A = \mathbb{C}[x] \) or \( \mathbb{C}[x]_{(x)} \) or \( \mathbb{C}(x) \). When do there exist \( g, h \in A \) such that \( f = gh \)?

First of all, we can take \( g = 1 \) and \( h = f \) or, more generally, \( g \) a unit in \( A \) and
Let us remark that this question is equivalent to the following: if \( \frac{A}{f} \) is an integral domain, is \( \frac{C[x]}{(f)C[x]} \) still an integral domain?

The answer to this question is no in general: set \( A := \mathbb{C}[x, y] \) and set \( f := x^2 - y^2(1 + y) \). Then \( f \) is irreducible as a polynomial since \( y^2(1 + y) \) is not a square in \( \mathbb{C}[x, y] \). But \( f = (x + y\sqrt{1 + y})(x - y\sqrt{1 + y}) \) where \( \sqrt{1 + y} \) is a formal power series such that \( \sqrt{1 + y^2} = 1 + y \). Thus \( f \) is not irreducible in \( \mathbb{C}[x, y] \) nor in \( \mathbb{C}\{x, y\} \) but it is irreducible in \( \mathbb{C}[x, y] \) or \( \mathbb{C}\{x, y\},(x,y) \).

In fact it is easy to see that \( x + y\sqrt{1 + y} \) and \( x - y\sqrt{1 + y} \) are power.series which are algebraic over \( \mathbb{C}[x, y] \), i.e. they are roots of polynomials with coefficients in \( \mathbb{C}[x, y] \). The set of such algebraic power series is a subring of \( \mathbb{C}[x, y] \) and it is denoted by \( \mathbb{C}\{x, y\} \). In general if \( x \) is a multivariable the ring of algebraic power series \( \mathbb{C}\{x\} \) is the following:

\[
\mathbb{C}\langle x \rangle := \{ f \in \mathbb{C}[x] \mid \exists P(z) \in \mathbb{C}[x][z], \ P(f) = 0 \}.
\]

It is not difficult to prove that the ring of algebraic power series is a subring of the ring of convergent power series and is a local ring. In 1969, M. Artin proved an analogue of Theorem 3.2.1 for the rings of algebraic power series [Ar69]. Thus if \( f \in \mathbb{C}\langle x \rangle \) (or \( \mathbb{C}\{x\} \)) is irreducible then it remains irreducible in \( \mathbb{C}\{x\} \), this is a consequence of Artin Approximation Theorem. From this theorem we can also deduce that if \( f \in C[x,I] \) (or \( C[x,I] \)), for some ideal \( I \), is irreducible, then it remains irreducible in \( \frac{C[x]}{I} \).

**Example 3.1.6.** — Let us strengthen the previous question. Let us assume that there exist \( \widehat{g}, \widehat{h} \in \mathbb{C}[x] \) such that \( f = \widehat{g}\widehat{h} \) with \( f \in A \) with \( A = \mathbb{C}\{x\} \) or \( \mathbb{C}\{x\} \). Then does there exist a unit \( \widehat{u} \in \mathbb{C}[x] \) such that \( \widehat{u}\widehat{g} \in A \) and \( \widehat{u}^{-1}\widehat{h} \in A \) ?

The answer to this question is positive if \( A = \mathbb{C}\langle x \rangle \) or \( \mathbb{C}\{x\} \), this is a non trivial corollary of Artin Approximation Theorem (see Corollary 3.4.4). But it is negative in general for \( \frac{C(x)}{I} \) or \( \frac{C\{x\}}{I} \) if \( I \) is an ideal. The following example is due to S. Izumi [Iz92]:

Set \( A := \frac{\mathbb{C}\{x,y,z\}}{(y^2 - xz^3)} \). Set \( \varphi(z) := \sum_{n=0}^{\infty} n!z^n \) (this is a divergent power series) and
set 
\( \hat{f} := x + y\hat{\varphi}(z), \quad \hat{g} := (x - y\hat{\varphi}(z))(1 - x\hat{\varphi}(z)^2)^{-1} \in \mathbb{C}[x, y, z]. \)

Then we can check that \( x^2 = \hat{f}\hat{g} \) modulo \( (y^2 - x^3) \). Now let us assume that there exists a unit \( \hat{u} \in \mathbb{C}[x, y, z] \) such that \( \hat{u}\hat{f} \in \mathbb{C}\{x, y, z\} \) modulo \( (y^2 - x^3) \). Thus \( P := \hat{u}\hat{f} - (y^2 - x^3)h \in \mathbb{C}\{x, y, z\} \) for some \( h \in \mathbb{C}[x, y, z] \). We can check easily that \( P(0, 0, 0) = 0 \) and \( \frac{\partial P}{\partial x}(0, 0, 0) = \hat{u}(0, 0, 0) \neq 0 \). Thus by the Implicit Function Theorem for analytic functions there exists \( \psi(y, z) \in \mathbb{C}\{y, z\} \), such that \( P(\psi(y, z), y, z) = 0 \) and \( \psi(0, 0) = 0 \). This yields 
\[ \psi(y, z) + y\hat{\varphi}(z) - (y^2 - \psi(y, z)^3)h(\psi(y, z), y, z)\hat{u}^{-1}(\psi(y, z), y, z) = 0. \]

By substituting 0 for \( y \) we obtain \( \psi(0, z) + \psi(0, z)^3\hat{k}(z) = 0 \) for some power series \( \hat{k}(z) \in \mathbb{C}[z] \). Since \( \psi(0, 0) = 0 \), this gives that \( \psi(0, z) = 0 \), thus \( \psi(y, z) = y\theta(y, z) \) with \( \theta(y, z) \in \mathbb{C}\{y, z\} \). Thus we obtain 
\[ \theta(y, z) + \hat{\varphi}(z) - (y - y^2\theta(y, z)^3)\hat{h}(\psi(y, z), y, z)\hat{u}^{-1}(\psi(y, z), y, z) = 0, \]
and by substituting 0 for \( y \), we see that \( \hat{\varphi}(z) = \theta(0, z) \in \mathbb{C}\{z\} \) which is a contradiction.

Thus \( x^2 = \hat{f}\hat{g} \) modulo \( (y^2 - x^3) \) but there is no unit \( \hat{u} \in \mathbb{C}[x, y, z] \) such that \( \hat{u}\hat{f} \in \mathbb{C}\{x, y, z\} \) modulo \( (y^2 - x^3) \).

**Example 3.1.7.** — A similar question is the following: if \( f \in A \) with \( A = \mathbb{C}[x], \mathbb{C}[x(x), C[x(x)] \) or \( C[x] \) and if there exist a non unit \( \hat{g} \in \mathbb{C}[x] \) and an integer \( m \in \mathbb{N} \) such that \( \hat{g}^m = f \), does there exist a non unit \( g \in A \) such that \( g^m = f \)?

A weaker question is the following: if \( \frac{A}{(f)} \) is reduced, is \( \frac{\mathbb{C}[x]}{(f)\mathbb{C}[x]} \) still reduced? Indeed, if \( \hat{g}^m = f \) for some non unit \( \hat{g} \) then \( \frac{\mathbb{C}[x]}{(f)\mathbb{C}[x]} \) is not reduced. Thus, if the answer to the second question is positive, then there exists a non unit \( g \in A \) and a unit \( u \in A \) such that \( u^kg^k = f \) for some integer \( k \).

As before, the answer to the first question is positive for \( A = \mathbb{C}(x) \) and \( A = \mathbb{C}\{x\} \) by Artin Approximation Theorem.

If \( A = \mathbb{C}[x] \) or \( \mathbb{C}[x(x) \), the answer to this question is negative. Indeed let us consider \( f = x^m + xn^{m+1} \). Then \( \hat{f} := g^m \) with \( \hat{g} := x \sqrt[n]{1 + x} \) but there is no \( g \in A \) such that \( g^m = f \).

Nevertheless, the answer to the second question is positive in the cases \( A = \mathbb{C}[x] \) or \( \mathbb{C}[x(x) \). This deep result is due to D. Rees (see [HuSw06] for instance).

**Example 3.1.8.** — Using the same notation as in Example 3.1.5 we can ask a stronger question: set \( A = \mathbb{C}(x) \) or \( \mathbb{C}\{x\} \) and let \( f \) be in \( A \). If there exist \( \hat{g} \) and \( \hat{h} \in \mathbb{C}[x] \), vanishing at 0, such that \( f = \hat{g}\hat{h} \) modulo a large power of the
ideal \((x)\), do there exist \(g\) and \(h\) in \(A\) such that \(f = gh\)? By example 3.1.5 there is no hope, if \(g\) and \(h\) exist, to expect that \(g\) and \(h \in \mathbb{C}[x]\).

We have the following theorem:

**Theorem 3.1.9 (Strong Artin Approximation Theorem)**

Let \(f(x,y)\) be a vector of polynomials over \(\mathbb{C}\) in two sets of variables \(x\) and \(y\). Then there exists a function \(\beta : \mathbb{N} \rightarrow \mathbb{N}\), such that for any integer \(c\) and any given approximate solution \(\tilde{y}(x)\) at order \(\beta(c)\),

\[
f(x, \tilde{y}(x)) \equiv 0 \mkern 1mu \text{modulo} \mkern 1mu (x)\beta(c),
\]

there exists an algebraic power series solution \(y(x)\),

\[
f(x, y(x)) = 0
\]

which coincides with \(\tilde{y}(x)\) up to degree \(c\),

\[
y(x) \equiv \tilde{y}(x) \mkern 1mu \text{modulo} \mkern 1mu (x)^c.
\]

In particular, if \(\tilde{y}h - f \equiv 0 \mkern 1mu \text{modulo} \mkern 1mu (x)\beta(1)\), where \(\beta\) is the function of the previous theorem for the polynomial \(y_1y_2 - f\), and if \(\tilde{y}(0) = \tilde{h}(0) = 0\), then there exist non units \(g\) and \(h \in \mathbb{C}(x)\) such that \(gh - f = 0\).

A natural question is: given \(f \in \mathbb{C}[x]\) how to compute \(\beta\) or, at least, \(\beta(1)\)? That is, up to what order do we have to check that the equation \(y_1y_2 - f = 0\) has an approximate solution in order to be sure that this equation has solutions? For instance, if \(f := x_1x_2 - x_3^d\) then \(f\) is irreducible but \(x_1x_2 - f \equiv 0 \mkern 1mu \text{modulo} \mkern 1mu (x)^d\) for any \(d \in \mathbb{N}\), so obviously \(\beta(1)\) really depends on \(f\).

In fact, in Theorem 3.1.9 M. Artin proved that \(\beta\) can be chosen independently of the degree of the components of the the vector \(f(x,y)\). But it is still an open problem to find effective bounds on \(\beta\) (see Section 3.3.4).

**Example 3.1.10 (Ideal Membership Problem).** — Set \(f_1, \ldots, f_r \in \mathbb{C}[x]\) where \(x = (x_1, \ldots, x_n)\). Let us denote by \(I\) the ideal of \(\mathbb{C}[x]\) generated by \(f_1, \ldots, f_r\). If \(g\) is a power series, how can we detect that \(g \in I\) or \(g \not\in I\)? Since a power series is determined by its coefficients, saying that \(g \in I\) will depend in general on an infinite number of conditions and it will not be possible to check that all these conditions are satisfied in finite time. Another problem is to find canonical representatives of power series modulo the ideal \(I\) that will help us to make computations in the quotient ring \(\mathbb{C}[x]/I\).

One way to solve these problems is the following. Let us consider the following order on \(\mathbb{N}^n\): for any \(\alpha, \beta \in \mathbb{N}^n\), we say that \(\alpha \leq \beta\) if \((|\alpha|, \alpha_1, \ldots, \alpha_n) \leq_{\text{lex}} (|\beta|, \beta_1, \ldots, \beta_n)\) where \(|\alpha| := \alpha_1 + \cdots + \alpha_n\) and \(\leq_{\text{lex}}\) is the lexicographic order. For instance

\[
(1, 1, 1) \leq (1, 2, 3) \leq (2, 2, 2) \leq (3, 2, 1) \leq (2, 2, 3).
\]
This order induces an order on the sets of monomials $x_1^{a_1} \ldots x_n^{a_n}$: we say that $x^\alpha \leq x^\beta$ if $\alpha \leq \beta$. Thus

$$x_1x_2x_3 \leq x_1x_2^2x_3 \leq x_1^2x_2x_3 \leq x_1x_2^2x_3^2.$$ 

If $f := \sum_{\alpha \in \mathbb{N}^n} f_\alpha x^\alpha \in \mathbb{C}[x]$, the initial exponent of $f$ with respect to the previous order is

$$\exp(f) := \min\{\alpha \in \mathbb{N}^n / f_\alpha \neq 0\} = \inf \text{Supp}(f)$$

where the support of $f$ is $\text{Supp}(f) := \{\alpha \in \mathbb{N}^n / f_\alpha \neq 0\}$. The initial term of $f$ is $f_{\exp(f)}$. This is the smallest non zero monomial in the Taylor expansion of $f$ with respect to the previous order.

If $I$ is an ideal of $\mathbb{C}[x]$, we define $\Gamma(I)$ to be the subset of $\mathbb{N}^n$ of all the initial exponents of elements of $I$. Since $I$ is an ideal, for any $\beta \in \mathbb{N}^n$ and any $f \in I$, $x^\beta f \in I$. This means that $\Gamma(I) + \mathbb{N}^n = \Gamma(I)$. Then we can prove that there exists a finite number of elements $g_1, \ldots, g_s \in I$ such that

$$\{\exp(g_1), \ldots, \exp(g_s)\} + \mathbb{N}^n = \Gamma(I).$$

Set

$$\Delta_1 := \exp(g_1) + \mathbb{N}^n \text{ and } \Delta_i = (\exp(g_i) + \mathbb{N}^n) \setminus \bigcup_{1 \leq j < i} \Delta_j, \text{ for } 2 \leq i \leq s.$$ 

Finally, set

$$\Delta_0 := \mathbb{N}^n \setminus \bigcup_{i=1}^s \Delta_i.$$ 

For instance, if $I$ is the ideal of $\mathbb{C}[x_1, x_2]$ generated by $g_1 := x_1x_2^3$ and $g_2 := x_1^2x_2^2$, we can check that

$$\Gamma(I) = \{(1, 3), (2, 2)\} + \mathbb{N}^2.$$
3.1. INTRODUCTION

Set \( g \in \mathbb{C}[x] \). Then, by Galligo-Grauert-Hironaka Division Theorem [Gal79], there exist unique power series \( q_1, \ldots, q_s, r \in \mathbb{C}[x] \) such that

\[
g = g_1q_1 + \cdots + g_sq_s + r
\]

\[
\exp(g_i) + \text{Supp}(q_i) \subset \Delta_i \text{ and } \text{Supp}(r) \subset \Delta_0.
\]

The uniqueness of the division comes from the fact the \( \Delta_i \)'s are disjoint subsets of \( \mathbb{N}^n \). The existence of such decomposition is proven through the division algorithm:

Set \( \alpha := \exp(g) \). Then there exists an integer \( i_1 \) such that \( \alpha \in \Delta_{i_1} \).

- If \( i_1 = 0 \), then set \( r^{(1)} := \text{in}(g) \) and \( q_i^{(1)} := 0 \) for all \( i \).
- If \( i_1 \geq 1 \), then set \( r^{(1)} := 0 \), \( q_i^{(1)} := 0 \) for \( i \neq i_1 \) and \( q_i^{(1)} := \frac{\text{in}(g)}{\exp(g_{i_1})} \).

Finally set \( g^{(1)} := g - \sum_{i=1}^s g_iq_i^{(1)} - r^{(1)} \). Thus we have \( \exp(g^{(1)}) > \exp(g) \). Then we replace \( g \) by \( g^{(1)} \) and the repeat the preceding process.

In this way we construct a sequence \( (g^{(k)})_k \) of power series such that, for any \( k \in \mathbb{N} \), \( \exp(g^{(k+1)}) > \exp(g^{(k)}) \) and \( g^{(k)} = g - \sum_{i=1}^s g_iq_i^{(k)} - r^{(k)} \) with

\[
\exp(g_i) + \text{Supp}(q_i^{(k)}) \subset \Delta_i \text{ and } \text{Supp}(r^{(k)}) \subset \Delta_0.
\]

At the limit \( k \to \infty \) we obtain the desired decomposition.

In particular since \( \{\exp(g_1), \ldots, \exp(g_s)\} + \mathbb{N}^n = \Gamma(I) \) we deduce from this that \( I \) is generated by \( g_1, \ldots, g_s \).

This algorithm means that for any \( g \in \mathbb{C}[x] \) there exists a unique power series \( r \) whose support is included in \( \Delta \) and such that \( g - r \in I \) and the division algorithm yields a way to obtain this representative \( r \).

Moreover, saying that \( g \not\in I \) is equivalent to \( r \neq 0 \) and this is equivalent to say that, for some integer \( k \), \( r^{(k)} \neq 0 \). But \( g \in I \) is equivalent to \( r = 0 \) which is equivalent to \( r^{(k)} = 0 \) for all \( k \in \mathbb{N} \). Thus applying the division algorithm, if for some integer \( k \), \( r^{(k)} \neq 0 \), then we can conclude that \( g \not\in I \). But this algorithm will not help us to determine if \( g \in I \) since we would have to make a infinite number of computations.

Now a natural question is, what happens if we replace \( \mathbb{C}[x] \) by \( A := \mathbb{C}(x) \) or \( \mathbb{C}\{x\} \)? Of course we can proceed with the division algorithm but we do not
know if \( q_1, \ldots, q_r \in A \). In fact by controlling the size of the coefficients of \( q_1^{(k)}, \ldots, q_r^{(k)}, r^{(k)} \) at each step of the division algorithm, we can prove that if \( g \in \mathbb{C}[x] \) then \( q_1, \ldots, q_r \) and \( r \) remain in \( \mathbb{C}[x] \) ([Hir64, Gra72, Gal79] and [JoPh00]). But if \( g \in \mathbb{C}(x) \) then it may happen that \( q_1, \ldots, q_r \) and \( r \) are not in \( \mathbb{C}(x) \) (see Example 3.5.4 of Section 3.5).

**Example 3.1.11 (Arcs Space and Jets Spaces).** — Let \( X \) be an affine algebraic subset of \( \mathbb{C}^m \), i.e. \( X \) is the zero locus of some polynomials in \( m \) variables: \( f_1, \ldots, f_r \in \mathbb{C}[y_1, \ldots, y_m] \). Let \( t \) be a single variable. For any integer \( n \), let us define \( X_n \) to be the set of vectors \( y(t) \) whose coordinates are polynomials of degree \( \leq n \) and such that \( f(y(t)) \equiv 0 \) modulo \( t^{n+1} \). The elements of \( X_n \) are called \( n \)-jets on \( X \).

If \( y_i(t) = y_{i,0} + y_{i,1}t + \cdots + y_{i,n}t^n \) and if we consider each \( y_{i,j} \) has one indeterminate, saying that \( f(y(t)) \in (t)^{n+1} \) is equivalent to the vanishing of \( r(n+1) \) polynomials equations involving the \( y_{i,j} \)'s. This shows that the jets spaces of \( X \) are algebraic sets.

For instance, if \( X \) is a cusp, i.e. the plane curve defined by \( X := \{ y_1^2 - y_2^3 = 0 \} \), then

\[
X_0 := \{(a_0, b_0) \in \mathbb{C}^2 \mid a_0^2 - b_0^3 = 0\} = X.
\]

We have

\[
X_1 = \{(a_0, a_1, b_0, b_1) \in \mathbb{C}^4 \mid (a_0 + a_1 t)^2 - (b_0 + b_1 t)^3 \equiv 0 \text{ modulo } t^2\}
\]

\[
= \{(a_0, a_1, b_0, b_1) \in \mathbb{C}^4 \mid a_0^2 - b_0^3 = 0 \text{ and } 2a_0a_1 - 3b_0^2b_1 = 0\}.
\]

The morphisms \( \mathbb{C}[t]/(t)^{n+1} \to \mathbb{C}[t]/(t)^{n+1} \), for \( k \geq n \), induce truncation maps \( \pi_n^k : X_k \to X_n \) by reducing \( k \)-jets modulo \( (t)^{n+1} \). In the example we are considering, the fibre of \( \pi_0^1 \) over the point \( (0, 1) \neq (0, 0) \) is the line in the \( (a_1, b_1) \)-plane whose equation is \( 2a_0a_1 - 3b_0^2b_1 = 0 \). This line is exactly the tangent space at \( X \) at the point \( (a_0, b_0) \). The tangent space at \( X \) in \( (0, 0) \) is the whole plane since this point is a singular point of the plane curve \( X \). This corresponds to the fact that the fibre of \( \pi_0^1 \) over \( (0, 0) \) is the whole plane.

On this example we show that \( X_1 \) is isomorphic to the tangent bundle of \( X \), which is a general fact.

We can easily see that \( X_2 \) is given by the following equations:

\[
\begin{align*}
    a_0^2 - b_0^3 &= 0 \\
    2a_0a_1 - 3b_0^2b_1 &= 0 \\
    a_1^2 + 2a_0a_2 - 3b_0^2b_2 &= 0
\end{align*}
\]

In particular, the fibre of \( \pi_0^2 \) over \( (0, 0) \) is the set of points of the form \( (0, 0, a_2, b_1, b_2) \) and the image of this fibre by \( \pi_1^2 \) is the line \( a_1 = 0 \). This shows that \( \pi_1^2 \) is not surjective.
But, we can show that above the smooth part of $X$, the maps $\pi^{n+1}$ are surjective and the fibres are isomorphic to $\mathbb{C}$.

The space of arcs on $X$, denoted by $X_\infty$, is the set of vectors $y(t)$ whose coordinates are formal power series satisfying $f(y(t)) = 0$. For such a general vector of formal power series $y(t)$, saying that $f(y(t)) = 0$ is equivalent to say that the coefficients of all the powers of $t$ in the Taylor expansion of $f(y(t))$ are equal to zero. This shows that $X_\infty$ may be defined by a countable number of equations in a countable number of variables. For instance, in the previous example, $X_\infty$ is the subset of $\mathbb{C}^N$ with coordinates $(a_0, a_1, a_2, ..., b_0, b_1, b_2, ...)$ defined by the infinite following equations:

$$
\begin{align*}
    a_0^2 - b_0^3 &= 0 \\
    2a_0a_1 - 3b_0^2b_1 &= 0 \\
    a_1^2 + 2a_0a_2 - 3b_0b_1^2 - 3b_0^2b_2 &= 0 \\
    \cdots &
\end{align*}
$$

The morphisms $\mathbb{C}[t] \to \mathbb{C}[t]/(t)^{n+1}$ induce truncations maps $\pi_n : X_\infty \to X_n$ by reducing arcs modulo $(t)^{n+1}$.

In general it is a difficult problem to compare $\pi_n(X_\infty)$ and $X_n$. It is not even clear if $\pi_n(X_\infty)$ is finitely defined. But we have the following theorem due to Greenberg which is a particular case of Theorem 3.1.9 in which $\beta$ is bounded by an affine function:

**Theorem 3.1.12 (Greenberg’s Theorem).** — [Gre66] Let $f(y)$ be a vector of polynomials in $m$ variables and let $t$ be a single variable. Then there exist two positive integers $a$ and $b$, such that for any polynomial solution $\overline{y}(t)$ modulo $(t)^{a+b}$, $f(\overline{y}(t)) \equiv 0$ modulo $(t)^{a+b+1}$, there exists a formal power series solution $\tilde{y}(t)$, $f(\tilde{y}(t)) = 0$ which coincides with $\overline{y}(t)$ up to degree $n + 1$, $\overline{y}(t) \equiv \tilde{y}(t)$ modulo $(t)^{n+1}$.

We can reinterpret this result as follows: let $X$ be the zero locus of $f$ and let $y(t)$ be a $(an + b)$-jet on $X$. Then the truncation of $y(t)$ modulo $(t)^{n+1}$ is the truncation of a formal power series solution of $f = 0$. Thus we have

$$
\pi_n(X_\infty) = \pi_n^{an+b}(X_{an+b}), \quad \forall n \in \mathbb{N}.
$$
A constructible subset of $\mathbb{C}^n$ is a set defined by the vanishing of some polynomials and the non-vanishing of other polynomials, i.e. a set of the form
\[
\{ x \in \mathbb{C}^n / f_1(x) = \cdots = f_r(x) = 0, g_1(x) \neq 0, \ldots, g_s(x) \neq 0 \}
\]
for some polynomials $f_i, g_j$. In particular algebraic sets are constructible sets.

Since a theorem of Chevalley asserts that the projection of an algebraic subset of $\mathbb{C}^{n+k}$ onto $\mathbb{C}^k$ is a constructible subset of $\mathbb{C}^n$, Theorem 3.1.12 asserts that $\pi_n(X_\infty)$ is a constructible subset of $\mathbb{C}^n$ since $X_{an+b}$ is an algebraic set. In particular $\pi_n(X_\infty)$ is finitely defined, i.e. it is defined by a finite number of data (see [GoLJ96] for an introduction to the study of these sets).

A difficult problem in singularity theory is to understand the behaviors of $X_n$ and $\pi_n(X_\infty)$ and to relate them to the geometry of $X$. One way to do this is to define the (motivic) measure of a constructible subset of $\mathbb{C}^n$, that is an additive map $\chi$ from the set of constructible sets to a commutative ring $R$, such that:

- $\chi(X) = \chi(Y)$ as soon as $X$ and $Y$ are isomorphic algebraic sets,
- $\chi(X\setminus U) + \chi(U) = \chi(X)$ as soon as $U$ is an open set of an algebraic set $X$,
- $\chi(X \times Y) = \chi(X)\chi(Y)$ for any algebraic sets $X$ and $Y$.

Then we are interested to understand the following formal power series:
\[
\sum_{n \in \mathbb{N}} \varphi(X_n)T^n \quad \text{and} \quad \sum_{n \in \mathbb{N}} \chi(\pi_n(X_\infty))T^n \in R[[T]].
\]

The reader may consult [DeLo99], [Lo00], [Ve06] for instance.

**Example 3.1.13.** — Let $f_1, \ldots, f_r \in k[x, y]$ where $k$ is an algebraically closed field and $x := (x_1, \ldots, x_n)$ and $y := (y_1, \ldots, y_m)$ are multivariables. Moreover we will assume here that $k$ is uncountable. As in the previous example let us define the following sets:
\[
X_l := \{ y(x) \in k[x]^m / f_i(x, y(x)) \in (x)^{l+1} \forall i \}.
\]

As we have done in the previous example, for any $l$ there exists an integer $N(l) \in \mathbb{N}$ such that $X_l \subset k^{N(l)}$. Moreover $X_l$ is an algebraic subset of $k^{N(l)}$ and the morphisms $\frac{k[x]}{(x)^{k+l}} \rightarrow \frac{k[x]}{(x)^{l}}$ for $k \leq l$ induce truncations maps $\pi^{X}_l : X_k \rightarrow X_l$ for any $k \geq l$.

By a theorem of Chevalley, for any $l \in \mathbb{N}$, the sequence $(\pi^{X}_l(X_k))_k$ is a decreasing sequence of constructible subsets of $X_l$. Thus the sequence $(\pi^{Y}_l(X_k))_k$ is a decreasing sequence of algebraic subsets of $X_l$, where $Y$ denotes the Zariski
closure of a subset \( Y \), i.e. the smallest algebraic set containing \( Y \). By Noetheri-
anity this sequence stabilizes: \( \pi_k^\ell(X_k) = \pi_{k'}^\ell(X_{k'}) \) for all \( k \) and \( k' \) large enough (say for any \( k, k' \geq k_l \)). Let us denote by \( F_l \) this algebraic set.

Let us assume that \( X_k \neq \emptyset \) for any \( k \in \mathbb{N} \). This implies that \( F_l \neq \emptyset \). Set \( C_{k,l} := \pi_k^l(X_k) \). It is a constructible set whose Zariski closure is \( F_l \) for any \( k \geq k_l \). Thus \( C_{k,l} \) has the form \( F_l \setminus V_k \) where \( V_k \) is an algebraic proper subset of \( F_l \), for any \( k \geq k_l \). Since \( \mathbb{k} \) is uncountable the set \( U_l := \bigcap_k C_{k,l} = \bigcap_k F_l \setminus V_k \) is not empty. By construction \( U_l \) is exactly the set of points of \( X_l \) that can be lifted to points of \( X_k \) for any \( k \geq l \). In particular \( \pi_k^l(U_k) = U_l \). If \( x_0 \in U_0 \) then \( x_0 \) may be lifted to \( U_1 \), i.e. there exists \( x_1 \in U_1 \) such that \( \pi_1^0(x_1) = x_0 \). By induction we may construct a sequence of points \( x_l \in U_l \) such that \( \pi_l^{l+1}(x_{l+1}) = x_l \) for any \( l \in \mathbb{N} \). At the limit we obtain a point \( x_\infty \) in \( X_\infty \), i.e. a power series \( y(x) \in \mathbb{k}[x]^{\mathbb{N}} \) solution of \( f(x,y) = 0 \).

We have proved here the following result similar to Theorem 3.1.9 if \( \mathbb{k} \) is an uncountable algebraically closed field and if \( f(x,y) = 0 \) has solutions modulo \( (x)^k \) for any \( k \in \mathbb{N} \), then there exists a power series solution \( y(x) \):

\[
f(x, y(x)) = 0.
\]

This kind of argument using asymptotic contructions (here the Noetherianity is the key point of the proof) may be nicely formalized using ultraproducts. Ultraproducts methods can be used to prove easily stronger results as Theorem 3.1.9 (See Part 3.3.3 and Proposition 3.3.24).

**Example 3.1.14 (Linearization of germ of diffeomorphism)**

Given \( f \in \mathbb{C}\{x\} \), \( x \) being a single variable, let us assume that \( f'(0) = \lambda \neq 0 \). Then \( f \) defines an analytic diffeomorphism from a neighborhood of \( 0 \) in \( \mathbb{C} \) onto a neighborhood of \( 0 \) in \( \mathbb{C} \) preserving the origin. The linearization problem, firstly investigated by C. L. Siegel, is the following: is \( f \) conjugated to its linear part? That is: does there exist \( g(x) \in \mathbb{C}\{x\} \), with \( g'(0) \neq 0 \), such that \( f(g(x)) = g(\lambda x) \) or \( g^{-1} \circ f \circ g(x) = \lambda x \) (in this case we say that \( f \) is analytically linearizable)?

This problem is difficult and the following cases may occur: \( f \) is not linearizable, \( f \) is formally linearizable but not analytically linearizable (i.e. \( g \) exists but \( g(x) \in \mathbb{C][x] \setminus \mathbb{C}\{x\} \), \( f \) is analytically linearizable (see [Ce91]).

Let us assume that \( f \) is formally linearizable, i.e. there exists \( \hat{g}(x) \in \mathbb{C}[x] \) such that \( f(\hat{g}(x)) - \hat{g}(\lambda x) = 0 \). By considering the Taylor expansion of \( \hat{g}(\lambda x) \):

\[
\hat{g}(\lambda x) = \hat{g}(y) + \sum_{n=1}^{\infty} \frac{(y - \lambda x)^n}{n!} f^{(n)}(y)
\]
we see that there exists \( \hat{h}(x, y) \in \mathbb{C}[x, y] \) such that \( \hat{g}(\lambda x) = \hat{g}(y) + (y - \lambda x)\hat{h}(x, y) \). Thus \( f \) is formally linearizable if and only if there exists \( \hat{h}(x, y) \in \mathbb{C}[x, y] \) such that

\[
f(\hat{g}(x)) - \hat{g}(y) + (y - \lambda x)\hat{h}(x, y) = 0.
\]

This former equation is equivalent to the existence of \( \hat{k}(y) \in \mathbb{C}[y] \) such that

\[
\begin{cases}
f(\hat{g}(x)) - \hat{k}(y) + (y - \lambda x)\hat{h}(x, y) = 0 \\
\hat{k}(y) - \hat{g}(y) = 0
\end{cases}
\]

Using the same trick as before (Taylor expansion), this is equivalent to the existence of \( \hat{l}(x, y, z) \in \mathbb{C}[x, y, z] \) such that

\[
\begin{cases}
f(\hat{g}(x)) - \hat{k}(y) + (y - \lambda x)\hat{h}(x, y) = 0 \\
\hat{k}(y) - \hat{g}(y) + (x - y)\hat{l}(x, y) = 0
\end{cases}
\]

Hence, we see that, if \( f \) is formally linearizable, there exists a formal solution \( (\hat{g}(x), \hat{k}(z), \hat{h}(x, y), \hat{l}(x, y, z)) \) of the system (1). Such a solution is called a solution with constraints. On the other hand, if the system (1) has a convergent solution \( (g(x), k(z), h(x, y), l(x, y, z)) \), then \( f \) is analytically linearizable.

We see that the problem of linearizing analytically \( f \) when \( f \) is formally linearizable is equivalent to find convergent power series solutions of the system (1) with constraints. Since it happens that \( f \) may be analytically linearizable but not formally linearizable, such a system (1) may have formal solutions with constraints but no analytic solutions with constraints.

In Section 3.5 we will give some results about the Artin Approximation Problem with constraints.

**Example 3.1.15.** — Another related problem is the following: if a differential equation with convergent power series coefficients has a formal power series solution, does it have convergent power series solutions? We can ask the same question by replacing "convergent" by "algebraic".

For instance let us consider the (divergent) formal power series \( \hat{y}(x) := \sum_{n=0}^{\infty} n!x^{n+1} \).

It is straightforward to check that it is a solution of the equation

\[
x^2y' - y + x = 0 \quad \text{(Euler Equation)}.
\]

On the other hand if \( \sum_n a_n x^n \) is a solution of the Euler Equation then the sequence \( (a_n)_n \) satisfies the following recursion:

\[
a_0 = 0, \quad a_1 = 1, \quad a_{n+1} = na_n \quad \forall n \geq 1.
\]
Thus $a_{n+1} = (n+1)!$ for any $n > 0$ and $\hat{y}(x)$ is the only solution of the Euler Equation. Hence we have an example of a differential equation with polynomials coefficients with a formal power series solution but without convergent power series solution. We will discuss in Section 3.5 how to relate this phenomenon to an Artin Approximation problem for polynomial equations with constraints (see Example 3.5.2).
Conventions We will assume that all the rings we consider are Noetherian commutative rings with unity. Ring morphisms $A \rightarrow B$ are assumed to take the unit element of $A$ into the unit element of $B$.

If $A$ is a local ring, then $m_A$ will denote its maximal ideal. For any $f \in A$, $f \neq 0$,

$$\text{ord}(f) := \max\{n \in \mathbb{N} \mid f \in m_A^n\}.$$  

If $A$ is an integral domain, $\text{Frac}(A)$ denotes its field of fractions.

If no other indication is given the letters $x$ and $y$ will always denote multivariables, $x := (x_1, ..., x_n)$ and $y := (y_1, ..., y_m)$, and $t$ will denote a single variable.

If $f(y)$ is a vector of polynomials with coefficients in a ring $A$,

$$f(y) := (f_1(y), ..., f_r(y)) \in A[y]^r,$$

if $\mathcal{I}$ is an ideal of $A$ and $\overline{y} \in A^m$, then $f(\overline{y}) \in \mathcal{I}$ (resp. $f(\overline{y}) = 0$) means $f_i(\overline{y}) \in \mathcal{I}$ (resp. $f_i(\overline{y}) = 0$) for $1 \leq i \leq r$.

3.2. Artin Approximation

In this first part we review the main results concerning the Artin Approximation Property. We give four results that are the most characteristic in the story: the classical Artin Approximation Theorem in the analytic case, its generalization by A. Płoski, a result of J. Denef and L. Lipschitz concerning rings with the Weierstrass Division Property and, finally, Popescu’s Approximation Theorem.

3.2.1. The analytic case. — In the analytic case, the first result is due to Michael Artin in 1968 \cite{Ar68}. His result asserts that the set of convergent solutions is dense in the set of formal solutions of a system of implicit analytic equations. This result is particularly useful, since if you have some analytic problem that you can express in a system of analytic equations, in order to find solutions of this problem you only need to find formal solutions and this may be done in general by an inductive process. Another way to use this result is the following: let us assume that you have some algebraic problem and that you are working over a ring of the form $A := k[[x]]$, where $x := (x_1, ..., x_n)$ and $k$ is a characteristic zero field. If the problem involves only a countable number of data (which is often the case in this context), since $\mathbb{C}$ is algebraically closed and the transcendence degree of $\mathbb{Q} \rightarrow \mathbb{C}$ is uncountable, you may assume that you work over $\mathbb{C}[x]$. Using Theorem \ref{3.2.1} you may, in some cases, reduce the problem to $A = \mathbb{C}\{x\}$. Then you can use powerful methods of complex analytic geometry to solve the problem. This kind of method is used, for instance, in the recent proof of the Nash Conjecture for algebraic surfaces (see Theorem A of \cite{FB12} and the crucial use of this theorem in \cite{FBPP12a}) or in the proof of the Abhyankar-Jung Theorem given in \cite{PaRo12}.
Let us mention that C. Chevalley had apparently proven this theorem some years before M. Artin but he did not publish it because he did not find applications of it [Ra].

3.2.1.1. Artin’s result. —

**Theorem 3.2.1.** — [Ar68] Let \( k \) be a valued field of characteristic zero and let \( f(x, y) \) be a vector of convergent power series in two sets of variables \( x \) and \( y \). Assume given a formal power series solution \( \hat{y}(x) \) vanishing at 0, \( f(x, \hat{y}(x)) = 0 \).

Then there exists, for any \( c \in \mathbb{N} \), a convergent power series solution \( y(x) \), \( f(x, \tilde{y}(x)) = 0 \) which coincides with \( \hat{y}(x) \) up to degree \( c \), \( \tilde{y}(x) \equiv \hat{y}(x) \) modulo \( (x)^c \).

**Remark 3.2.2.** — This theorem has been conjectured by S. Lang in [Lan54] (last paragraph p. 372) when \( k = \mathbb{C} \).

**Remark 3.2.3.** — The ideal \( (x) \) defines a topology on \( k^n \times k^m \) called the Krull topology induced by the following norm: \( |a(x)| := e^{-\text{ord}(a(x))} \). In this case small elements of \( k[x] \) are elements of high order. Thus Theorem 3.2.1 asserts that the set of solutions in \( k^n \times k^m \) of \( f(x, y) = 0 \) is dense in the set of solutions in \( k^n \times k^m \) of \( f(x, y) = 0 \) for the Krull topology.

**Proof of Theorem 3.2.1.** — Let us first give the main ideas of the proof. The proof is done by induction on \( n \), the case \( n = 0 \) being obvious.

The first step is to reduce the problem to the case the ideal \( I \) generated by \( f_1, \ldots, f_r \) is a prime ideal by adding to \( I \) all the elements \( g(x, y) \) such that \( g(x, \hat{y}(x)) = 0 \). Let us denote by \( X \) the analytic set defined by \( I \).

The next step is to reduce to the case \( X \) is complete intersection, this means that \( I \) is generated by \( r \) elements where \( r \) is equal to the codimension of \( X \) in \( k^n \times k^m \).

After these reductions, the proper proof starts. The key ingredient is a suitable minor \( \delta \) of the Jacobian matrix \( \left( \frac{\partial f}{\partial y} \right) \) of \( f \), namely one which is not identically zero on \( X \). The existence of such a minor is ensured by the Jacobian Criterion: at a smooth point of \( X \), the rank of the Jacobian matrix is the codimension of \( X \) at this point. Since the set of smooth points is dense, the assertion follows. We denote by \( \hat{\delta}(x) := \delta(x, \hat{y}(x)) \) the evaluation of \( \delta \) at our given formal solution. Then, the idea is the following: instead of trying to solve \( f(x, y) = 0 \) with a convergent solution, we aim at finding a convergent power series vector \( \tilde{y}(x) \) such that \( \delta^2(x, \tilde{y}(x)) \) divides \( f(x, \tilde{y}(x)) \). Since \( f(x, \hat{y}(x)) = 0 \), then \( \delta^2(x, \tilde{y}(x)) \) already divides \( f(x, \tilde{y}(x)) \), we will reformulate the statement.


Let \( I \) be the ideal of \( k\{x,y\} \) generated by \( f_1(x,y), \ldots, f_r(x,y) \). Let \( \varphi \) be the \( k\{x,y\} \)-morphism \( k\{x,y\} \rightarrow k[[x]] \) sending \( y_i \) onto \( \hat{y}_i(x) \). Then \( \text{Ker}(\varphi) \) is a prime ideal containing \( I \) and if the theorem is true for generators of \( \text{Ker}(\varphi) \) then it is true for \( f_1, \ldots, f_r \). Thus we can assume that \( I = \text{Ker}(\varphi) \). The local ring \( k\{x,y\}_I \) is regular by a theorem of Serre (see Theorem 19.3 [Mat80]). Set \( h := \text{height}(I) \). Thus, from the Jacobian Criterion, there exists a \( h \times h \) minor of the Jacobian matrix \( \frac{\partial(f_1, \ldots, f_r)}{\partial(y)} \), denoted by \( \delta(x,y) \), such that \( \delta \notin I = \text{Ker}(\varphi) \). In particular we have \( \delta(x,\hat{y}(x)) \neq 0 \).

By considering the partial derivative of \( f_i(x,\hat{y}(x)) = 0 \) with respect to \( x_j \) we get

\[
\frac{\partial f_i}{\partial x_j}(x,\hat{y}(x)) = - \sum_{k=1}^r \frac{\partial \hat{y}_k(x)}{\partial x_j} \frac{\partial f_i}{\partial y_k}(x,\hat{y}(x)).
\]

Thus there exists a \( h \times h \) minor of the Jacobian matrix \( \frac{\partial(f_1, \ldots, f_r)}{\partial(y)} \), still denoted by \( \delta(x,y) \), such that \( \delta(x,\hat{y}(x)) \neq 0 \). In particular \( \delta \notin I \). From now on we will assume that \( \delta \) is the determinant of \( \frac{\partial(f_1, \ldots, f_k)}{\partial(y_1, \ldots, y_k)} \).

If we denote \( J := (f_1, \ldots, f_k) \), then \( \text{ht}(Jk\{x,y\}_I) \leq h \). On the other hand we have \( \text{ht}(Jk\{x,y\}_I) \geq \text{rk}(\frac{\partial(f_1, \ldots, f_k)}{\partial(y_1, \ldots, y_k)}) \mod I \), and \( h \leq \text{rk}(\frac{\partial(f_1, \ldots, f_k)}{\partial(y_1, \ldots, y_k)}) \mod I \).
since $\delta(x, \hat{g}(x)) \neq 0$. Thus $\text{ht}(Jk\{x, y\}) = h$ and $\sqrt{Jk\{x, y\}_I} = Ik\{x, y\}_I$.

This means that there exists $q \in k\{x, y\}$, $q \notin I$, and $e \in \mathbb{N}$ such that $qf^e_i \in J$ for $h + 1 \leq i \leq m$. In particular $q(x, \hat{g}(x)) \neq 0$. We will use this fact later.

Then we will use the following lemma with $g := \delta^2$.

**Lemma 3.2.4.** — Let us assume that Theorem 3.2.1 is true for an integer $n - 1$. Let $g(x, y)$ be a convergent power series and let $f(x, y)$ be a vector of convergent power series.

Let $\hat{g}(x)$ be in $(x)k[[x]]^m$ such that $g(x, \hat{g}(x)) \neq 0$ and $f(x, \hat{g}(x)) = 0 \mod. g(x, \hat{g}(x))$.

Let $c$ be an integer. Then there exists $\overline{g}(x) \in (x)k\{x\}^m$ such that $f(x, \overline{g}(x)) = 0 \mod. g(x, \overline{g}(x))$ and $\overline{g}(x) - \hat{g}(x) \in (x)^c$.

**Proof of Lemma 3.2.4.** — If $g(x, \hat{g}(x))$ is invertible, the result is obvious (just take for $\hat{g}(x)$ any truncation of $\hat{g}(x)$). Thus let us assume that $g(x, \hat{g}(x))$ is not invertible. By making a linear change of variables we may assume that $g(x, \hat{g}(x))$ is regular with respect to $x_n$ and by Weierstrass Preparation Theorem $g(x, \hat{g}(x)) = \hat{a}(x) \times \text{unit}$ where

$$\hat{a}(x) := x_n^d + \hat{a}_1(x') x_n^{d-1} + \cdots + \hat{a}_d(x')$$

where $x' := (x_1, \ldots, x_{n-1})$ and $a_i(x') \in (x')k[[x']]$, $1 \leq i \leq d$. Let us perform the Weierstrass division of $\hat{g}(x)$ by $\hat{a}(x)$:

$$\hat{g}_i(x) = \hat{a}(x) \hat{w}_i(x) + \sum_{j=0}^{d-1} \hat{g}_{i,j}(x') x_n^j$$

for $1 \leq i \leq m$. Let us denote

$$\hat{g}^*(x)_i := \sum_{j=0}^{d-1} \hat{g}_{i,j}(x') x_n^j, \quad 1 \leq i \leq m.$$

Then $g(x, \hat{g}(x)) = g(x, \hat{g}^*(x)) \mod. \hat{a}(x)$ and $f_k(x, \hat{g}(x)) = f_k(x, \hat{g}^*(x)) \mod. \hat{a}(x)$ for $1 \leq k \leq r$.

Let $y_{i,j}$, $1 \leq i \leq m$, $1 \leq j \leq d - 1$, be new variables. Let us denote $y^*_i := \sum_{j=1}^{d-1} y_{i,j} x_n^j, \quad 1 \leq i \leq m$. Let us denote the polynomial

$$A(a_i, x_n) := x_n^d + a_1 x_n^{d-1} + \cdots + a_d \in k[x_n, a_1, \ldots, a_d]$$

where $a_1, \ldots, a_d$ are new variables. Let us perform the Weierstrass division of $g(x, y^*)$ and $f_i(x, y^*)$ by $A$:

$$g(x, y^*) = A.Q + \sum_{l=1}^{d-1} G lx_n^j$$
\[ f_k(x, y^*) = A(Q_k + \sum_{l=1}^{d-1} F_{k,l}x_n^l), \quad 1 \leq k \leq r \]

where \( Q, Q_k \in \mathbb{k}\{x, y_{i,j}, a_p\} \) and \( G_l, F_{k,l} \in \mathbb{k}\{x', y_{i,j}, a_p\} \). Then we have

\[ g(x, \tilde{y}^*(x)) = \sum_{l=1}^{d-1} G_l(x', \tilde{y}_{i,j}(x'), \tilde{a}_p(x'))x_n^l \mod. (\tilde{a}(x)) \]

\[ f_k(x, \tilde{y}^*(x)) = \sum_{l=1}^{d-1} F_{k,l}(x', \tilde{y}_{i,j}(x'), \tilde{a}_p(x'))x_n^l \mod. (\tilde{a}(x)), \quad 1 \leq k \leq r. \]

This proves that \( G_l(x', \tilde{y}_{i,j}(x'), \tilde{a}_p(x')) = 0 \) and \( F_{k,l}(x', \tilde{y}_{i,j}(x'), \tilde{a}_p(x')) = 0 \) for all \( k \) and \( l \). By the inductive hypothesis, there exists \( \overline{y}_{i,j}(x') \in \mathbb{k}\{x'\} \) and \( \overline{a}_p(x') \in \mathbb{k}\{x'\} \) for all \( i, j \) and \( s \), such that \( G_l(x', \overline{y}_{i,j}(x'), \overline{a}_p(x')) = 0 \) and \( F_{k,l}(x', \overline{y}_{i,j}(x'), \overline{a}_p(x')) = 0 \) for all \( k \) and \( l \) and \( \overline{y}_{i,j}(x') - \tilde{y}_{i,j}(x'), \overline{a}_p(x') - \tilde{a}_p(x') \in (x')^c \) for all \( i, j \) and \( p \) (Formally in order to apply the induction hypothesis we should have \( \tilde{y}_{i,j}(0) = 0 \) and \( \tilde{a}_p(0) = 0 \) which is not necessarily the case here.

We can remove the problem by replacing \( \tilde{y}_{i,j}(x') \) and \( \tilde{a}_p(x') \) by \( \overline{y}_{i,j}(x') - \tilde{y}_{i,j}(0) \) and \( \overline{a}_p(x') - \tilde{a}_p(0) \), and \( G_l(x', y_{i,j}, a_p) \) by \( G'(x', y_{i,j} + \tilde{y}_{i,j}(0), a_p + \tilde{a}_p(0)) \) - idem for \( F_{k,l} \).

Let us denote

\[ \overline{a}(x) := x_n^d + \overline{a}_1(x)x_n^{d-1} + \cdots + \overline{a}_d(x') \]

\[ \overline{y}_{i}(x) := \overline{a}(x)\overline{w}_{i}(x) + \sum_{j=0}^{d-1} \overline{y}_{i,j}(x')x_n^j \]

for some \( \overline{w}_{i}(x) \in \mathbb{k}\{x\} \) such that \( \overline{w}_{i}(x) - \tilde{w}_{i}(x) \in (x)^c \) for all \( i \). It is straightforward to check that \( f_i(x, \overline{y}(x)) = 0 \mod. g(x, \overline{y}(x)) \) for \( 1 \leq i \leq r \) and \( \overline{y}_{j}(x) - \tilde{y}_{j}(x) \in (x)^c \) for \( 1 \leq j \leq m \).

We can apply this lemma to \( g(x, y) := \delta^2(x, y) \) with \( c' := c + d + 1 \) and \( d := \text{ord}(\delta^2(x, \tilde{y}(x))) \). Thus we may assume that there is \( \overline{y}_{i}(x) \in \mathbb{k}\{x\}, 1 \leq i \leq m, \) such that \( f(x, \overline{y}) \in \delta^2(x, \overline{y}) \) and \( \overline{y}_{i}(x) - \tilde{y}_{i}(x) \in (x)^{c'+d+1}, 1 \leq i \leq m. \) Since \( \text{ord}(\delta^2(x, \overline{y})) = d \), then we have \( f(x, \overline{y}) \in \delta^2(x, \overline{y})(x)^c \). Then we use the following generalization of the Implicit Function Theorem to show that there exists \( \tilde{y}(x) \in \mathbb{k}\{x\}^m \) with \( \tilde{y}(0) = 0 \) such that \( \overline{y}_{j}(x) - \tilde{y}_{j}(x) \in (x)^c, 1 \leq j \leq m, \) and and \( f_i(x, \tilde{y}(x)) = 0 \) for \( 1 \leq i \leq h. \)

**Theorem 3.2.5** (Tougeron’s Implicit Function Theorem)
Let \( f(x, y) \) be a vector of \( \mathbb{k}\{x, y\}^h \) with \( m \geq h \), and let \( \delta(x, y) \) be a \( h \times h \) minor of the Jacobian matrix \( \frac{\partial (f_1, \ldots, f_h)}{\partial (y_1, \ldots, y_m)} \). Let us assume that there exists \( g(x) \in \mathbb{k}\{x\}^m \) such that

\[
\frac{\partial (f_1, \ldots, f_h)}{\partial (y_1, \ldots, y_m)}(x, y) = 0 \quad \text{for all } 1 \leq i \leq h
\]

and for some \( c \in \mathbb{N} \). Then there exists \( \tilde{y}(x) \in \mathbb{k}\{x\}^m \) such that

\[
f_i(x, \tilde{y}(x)) = 0 \quad \text{for all } 1 \leq i \leq h
\]

Moreover \( \tilde{y}(x) \) is unique if we impose \( \tilde{y}_j(x) = y_j(x) \) for \( h < j \leq m \).

If \( c > \text{ord}(q(x, \tilde{y}(x))) \), then \( q(x, \tilde{y}(x)) \neq 0 \). Since \( q f_i^e \in J \) for \( h + 1 \leq i \leq r \), this proves that \( f_i(x, \tilde{y}(x)) = 0 \) for all \( i \).

**Proof of Theorem 3.2.4** — We may assume that \( \delta \) is the first \( r \times r \) minor of the Jacobian matrix. If we add the equations \( f_{h+1} := y_{h+1} - \tilde{y}_{h+1}(x) = 0, \ldots, f_m := y_m - \tilde{y}_m(x) = 0 \), we may assume that \( m = h \) and \( \delta \) is the determinant of the Jacobian matrix \( J(x, y) := \frac{\partial (f_1, \ldots, f_h)}{\partial (y_1, y_2)} \). We have

\[
f(x, y(x) + \delta(x, y(x))z) = f(x, y(x)) + \delta(x, y)z J(x, y(x)) + \delta(x, y(x)) \delta(h(x, y(x)), z)
\]

where \( z := (z_1, \ldots, z_m) \) and \( H(x, y(x), z) \in \mathbb{k}\{x, y(x), z\}^m \) is of order at least \( 2 \) in \( z \). Let us denote by \( J'(x, y(x)) \) the comatrix of \( J(x, y(x)) \). Let \( \varepsilon(x) \) be in \( (x)^r \mathbb{k}\{x\}^r \) such that \( f(x, y(x)) = \delta^2(x, y(x)) \varepsilon(x) \). Then we have

\[
f(x, y(x) + \delta(x, y(x))z) = \\
= \delta(x, y(x)) \left( \varepsilon(x) J'(x, y(x)) + z + H(x, y(x), z) J'(x, y(x)) \right) J(x, y(x)).
\]

Let us denote

\[
g(x, z) := \varepsilon(x) J'(x, y(x)) + z + H(x, y(x), z) J'(x, y(x)).
\]

Then \( g(0, 0) = 0 \) and the matrix \( \frac{\partial g(x, z)}{\partial z}(0, 0) \) is the identity matrix. Thus, by the Implicit Function Theorem, there exists a unique \( z(x) \in \mathbb{k}\{x\}^m \) such that \( f(x, y(x) + \delta(x, y(x))z(x)) = 0 \). This proves the theorem.

**Remark 3.2.6.** — We can do the following remarks about the proof of Theorem 3.2.1

i) In the case \( n = 1 \) i.e. \( x \) is a single variable, set \( e := \text{ord}(\delta(x, \tilde{y}(x))) \). If \( \overline{y}(x) \in \mathbb{k}\{x\}^m \) satisfies \( \tilde{y}(x) - \overline{y}(x) \in (x)^{2e+c} \), then we have

\[
\text{ord}(f(x, \overline{y}(x))) \geq 2e + c
\]

and

\[
\delta(x, \overline{y}(x)) = \delta(x, \tilde{y}(x)) \text{ mod. } (x)^{2e+c},
\]
thus \( \text{ord}(\delta(x, y(x))) = \text{ord}(\delta(x, \hat{y}(x))) = e. \) Hence we have automatically \( f(x, y(x)) \in (\delta(x, y(x)))^2(x)^e \) since \( k\{x\} \) is a discrete valuation ring (i.e. if \( \text{ord}(a(x)) \leq \text{ord}(b(x)) \) then \( a(x) \) divides \( b(y) \) in \( k\{x\} \)).

Thus Lemma \( 3.2.4 \) is not necessary in this case and the proof is quite simple. This fact will be general: approximation results will be easier to obtain, and sometimes stronger, in discrete valuation rings than in more general rings.

ii) In fact, we did not use that \( k \) is a field of characteristic zero, we just need \( k \) to be a perfect field in order to use the Jacobian Criterion. But the use of the Jacobian Criterion is more delicate for non perfect fields. This also will be general: approximation results will be more difficult to prove in positive characteristic. For instance M. André proved Theorem \( 3.2.1 \) in the case \( k \) is a complete field of positive characteristic and replace the use of the Jacobian Criterion by the homology of commutative algebras \( \text{[An75]} \).

iii) For \( n \geq 2 \), the proof of Theorem \( \text{[Ar68]} \) uses an induction on \( n \). In order to do it we use the Weierstrass Preparation Theorem. But to apply the Weierstrass Preparation Theorem we need to do a linear change of coordinates in \( k\{x\} \), in order to transform \( g(x, \hat{y}(x)) \) into a power series \( h(x) \) such that \( h(0, ..., 0, x_n) \neq 0 \). Then the proof does not adapt to prove similar results in the case of constraints: for instance if \( \hat{y}_1(x) \) depends only on \( x_1 \) and \( \hat{y}_2(x) \) depends only on \( x_2 \), can we find a convergent solution such that \( \hat{y}_1(x) \) depends only on \( x_1 \), and \( \hat{y}_2(x) \) depends only on \( x_2 \)? Moreover, even if we can use a linear change of coordinates without modifying the constrains, the use of the Tougeron’s Implicit Function Theorem may remove the constrains. We will discuss these problems in Section \( 3.5 \).

**Corollary 3.2.7.** — Let \( k \) be a valued field of characteristic zero and let \( I \) be an ideal of \( k\{x\} \). If \( f(y) \in \left( \frac{k\{x, y\}}{IK(x, y)} \right)^r \), let \( \hat{y} \in \left( \frac{k\{x\}}{IK(x)} \right)^m \) be a solution of \( f = 0 \) such that \( \hat{y} \equiv 0 \) modulo \( I + (x) \). Then there exists a solution of \( f = 0 \) in \( \frac{k\{x\}^m}{IK(x)} \) denoted by \( \hat{y} \) such that \( \hat{y} \equiv 0 \) modulo \( I + (x) \) and \( \hat{y} - \hat{y} \in (x)^e \frac{k\{x\}^m}{IK(x)} \).

**Proof.** — Set \( F_i(x, y) \in k\{x, y\} \) such that \( F_i(x, y) = f_i(y) \) mod. \( I \) for \( 1 \leq i \leq r \). Let \( a_1, ..., a_s \in k\{x\} \) be generators of \( I \). Set \( \hat{w}(x) \in k[x]^m \) such that \( \hat{w}_j(x) \equiv \hat{y}_j \) mod. \( I \) for \( 1 \leq j \leq m \). Since \( f_i(\hat{y}) = 0 \) then there exists \( \hat{z}_{i,k}(x) \in k[x] \), \( 1 \leq i \leq r \) and \( 1 \leq k \leq s \), such that

\[
F_i(x, \hat{w}(x)) + a_1 \hat{z}_{i,1}(x) + \cdots + a_s \hat{z}_{i,s}(x) = 0 \quad \forall i.
\]

After Theorem \( 3.2.1 \) there exist \( \hat{w}_j(x), \hat{z}_{i,k}(x) \in k\{x\} \) such that

\[
F_i(x, \hat{w}(x)) + a_1 \hat{z}_{i,1}(x) + \cdots + a_s \hat{z}_{i,s}(x) = 0 \quad \forall i.
\]
and \( \tilde{w}_j(x) - \tilde{w}_j(x) \in (x)^c \) for \( 1 \leq j \leq m \). Then the images of the \( \tilde{w}_j(x) \)'s in \( \frac{k[x]}{I} \) satisfy the conclusion of the corollary.

3.2.1.2. Płoski’s result. — A few years after M. Artin’s result, A. Płoski strengthened Theorem 3.2.1 by a careful analysis of the proof. His result yields an analytic parametrization of a piece of the set of solutions of \( f = 0 \) such that the formal solution \( \tilde{w}(x) \) is a formal point of this parametrization.

**Theorem 3.2.8.** — [Pł74] Let \( k \) be a valued field of characteristic zero and let \( f(x, y) \) be a vector of power series in two in \( k\{x, y\}^r \). Let \( \tilde{y}(x) \) be a formal power series solution such that \( \tilde{y}(0) = 0 \), \( f(x, \tilde{y}(x)) = 0 \).

Then there exists a convergent power series solution \( y(x, z) \in k\{x, z\}^m \), where \( z = (z_1, ..., z_s) \) are new variables, \( f(x, y(x, z)) = 0 \), and a vector of formal power series \( \tilde{z}(x) \in k[x]^s \) with \( \tilde{z}(0) = 0 \) such that

\[
\tilde{y}(x) = y(x, \tilde{z}(x)).
\]

This result obviously implies Theorem 3.2.1 since we can choose convergent power series \( \tilde{z}_1(x), ..., \tilde{z}_s(x) \in k\{x\}^s \) such that \( \tilde{z}_j(x) = \tilde{w}_j(x) \in (x)^c \) for \( 1 \leq j \leq s \). Then, by denoting \( \tilde{y}(x) := y(x, \tilde{z}(x)) \), we get the conclusion of Theorem 3.2.1.

**Remark 3.2.9.** — Let us remark that this result remains true if we replace \( k\{x\} \) by a quotient \( \frac{k(x)}{I} \) as in Corollary 3.2.7.

**Remark 3.2.10.** — Let \( I \) be the ideal generated by \( f_1, ..., f_r \). The formal solution \( \tilde{y}(x) \) of \( f = 0 \) induces a \( k\{x\} \)-morphism \( k\{x, y\} \rightarrow k[x] \) defined by the substitution of \( y(x) \) for \( y \). Then \( I \) is included in the kernel of this morphism thus, by the universal property of the quotient ring, this morphism induces a \( k\{x\} \)-morphism \( \psi : \frac{k[x, y]}{I} \rightarrow k[x] \). On the other hand, any \( k\{x\} \)-morphism \( \psi : \frac{k[x, y]}{I} \rightarrow k[x] \) is clearly defined by substituting for \( y \) a formal power series \( \tilde{y}(x) \) such that \( f(x, \tilde{y}(x)) = 0 \).

Thus we can reformulate Theorem 3.2.8 as follows: Let \( \psi : \frac{k[x, y]}{I} \rightarrow k[x] \) be the \( k\{x\} \)-morphism defined by the formal power series solution \( \tilde{y}(x) \). Then there exist an analytic \( k\{x\} \)-algebra \( D := k\{x, z\} \) and \( k\{x\} \)-morphisms \( C \rightarrow D \) (defined via the convergent power series solution \( y(x, z) \) of \( f = 0 \)) and \( D \rightarrow k[x] \) (defined by substituting \( \tilde{z}(x) \) for \( z \)) such that the following diagram commutes:
3.2.2. Artin Approximation and Weierstrass Division Theorem. —

The proof of Theorem 3.2.1 uses essentially only two results: the Weierstrass Division Theorem and the Implicit Function Theorem. In particular it is straightforward to check that the proof of Theorem 3.2.1 remains true if we replace \( k\{x,y\} \) by \( k(x,y) \), the ring of algebraic power series in \( x \) and \( y \), since this ring satisfies the Weierstrass Division Theorem (cf. [La67], see Section 3.6) and the Implicit Function Theorem. In [Ar69], M. Artin gives a version of Theorem 3.2.1 in the case of polynomials equations over a field or an excellent discrete valuation ring \( k \), and proves that formal solutions of such equations can be approximated by solutions in the Henselization of the ring of polynomials over \( k \), i.e. in a localization of a finite extension of the ring of polynomials over \( k \). The proof, when \( k \) is an excellent discrete valuation ring, uses Néron \( p \)-desingularization [Né64] (see Section 3.2.3 for a statement of Néron \( p \)-desingularization). This result is very important since it allows to reduce some algebraic problems over complete local ring to local rings which are localization of finitely generated rings over a field or a discrete valuation ring.

For instance, this idea, along with an idea of C. Peskine and L. Szpiro, was used by M. Hochster to reduce problems over complete local rings in characteristic zero to the same problems in positive characteristic. The idea is the following: let us assume that some statement \((T)\) is true in positive characteristic (where you can use the Frobenius map to prove it for instance) and let us assume that there exists an example showing that \((T)\) is not true in characteristic zero. In some cases we can use Artin Approximation Theorem to show the existence of a counterexample to \((T)\) in the Henselization at a prime ideal of a finitely generated algebra over a field of characteristic zero. Since the Henselization is the direct limit of étale extensions, we can show the existence of a counterexample to \((T)\) in a local ring \( A \) which is the localization of a finitely generated algebra over a field of characteristic zero \( k \). If the example involves only a finite number of data in \( A \), then we may lift this counterexample in a ring which is the localization of a finitely generated ring over \( \mathbb{Q} \), and even over \( \mathbb{Z}/p\mathbb{Z} \) for all but finitely many

\[
\begin{align*}
\mathbb{k}\{x\} & \xrightarrow{\varphi} \mathbb{k}[x] \\
\mathbb{k}\{x,y\} & \xrightarrow{\psi} D := \mathbb{k}\{x, z\}
\end{align*}
\]

We will use and generalize this formulation later (see Theorem 3.2.16).
primes \( p \) by reducing the problem modulo \( p \) (in fact for \( p \neq p_i \) for \( 1 \leq i \leq s \)).

This idea was used to prove important results about Intersection Conjectures [PeSz73], big Cohen-Macaulay modules [HoRo74], Homological Conjectures [Ho75].

J Denef and L. Lipschitz axiomatized the properties a ring needs to satisfy in order to adapt the proof the main theorem of [Ar69] due M. Artin. They called such families of rings Weierstrass Systems. There are two reasons for introducing such rings: the first one is the proof of Theorem 3.5.16 and the second one is their use in proofs of Strong Artin Approximation results. Independently H. Kurke, G. Pfister, D. Popescu, M. Roczen and T. Mostowski (cf. [KPPRM78]) introduced the notion of Weierstrass category which is very similar (see [KuPf82] for a connection between these two notions).

**Definition 3.2.11.** — [DeLi80] Let \( k \) be a field or a discrete valuation ring of maximal ideal \( \mathfrak{p} \). By a **Weierstrass System of local \( k \)-algebras**, or a **W-system** over \( k \), we mean a family of \( k \)-algebras \( k[[x_1, \ldots, x_n]] \), \( n \in \mathbb{N} \) such that:

i) For \( n = 0 \), the \( k \)-algebra is \( k \).

For any \( n \geq 1 \), \( k[x_1, \ldots, x_n]/(p, x_1, \ldots, x_n) \subset k[[x_1, \ldots, x_n]] \subset k[[x_1, \ldots, x_n]] \) and \( k[[x_1, \ldots, x_{n+m}]] \cap k[[x_1, \ldots, x_n]] = k[[x_1, \ldots, x_n]] \) for \( m \in \mathbb{N} \). For any permutation \( \sigma \) of \( \{1, \ldots, n\} \) if \( f \in k[[x_1, \ldots, x_n]] \), then \( f(x_{\sigma(1)}, \ldots, x_{\sigma(n)}) \in k[[x_1, \ldots, x_n]] \).

ii) Any element of \( k[[x]], x = (x_1, \ldots, x_n) \), which is a unit in \( k[[x]] \), is a unit in \( k[[x]] \).

iii) If \( f \in k[[x]] \) and \( \mathfrak{p} \) divides \( f \) in \( k[[x]] \) then \( \mathfrak{p} \) divides \( f \) in \( k[[x]] \).

iv) Let \( f \in (p, x)k[[x]] \) such that \( f \neq 0 \). Suppose that \( f \in (p, x_1, \ldots, x_{n-1}, x_n^s) \) but \( f \notin (p, x_1, \ldots, x_{n-1}, x_n^{s-1}) \). Then for any \( g \in k[[x]] \) there exist a unique \( q \in k[[x]] \) and a unique \( r \in k[[x_1, \ldots, x_{n-1}]] \) with \( \deg x_r < d \) such that \( g = qf + r \).

v) (if \( \text{char}(k) > 0 \)) If \( \bar{g} \in (p, x_1, \ldots, x_n)k[[x_1, \ldots, x_n]] \) and \( f \in k[[y_1, \ldots, y_m]] \) such that \( f \neq 0 \) and \( f(\bar{g}) = 0 \), then there exists \( g \in k[[y]] \) irreducible in \( k[[y]] \) such that \( g(\bar{y}) = 0 \) and such that there does not exist any unit \( u(y) \in k[[y]] \) with \( u(y)g(y) = \sum_{a \in \mathbb{N}^n} a_\alpha y^{\alpha} \) (\( a_\alpha \in k \)).

vi) (if \( \text{char}(k/p) \neq 0 \)) Let \( (k/p)[[x]] \) be the image of \( k[[x]] \) under the projection \( k[[x]] \twoheadrightarrow (k/p)[[x]] \). Then \( (k/p)[[x]] \) satisfies \( v \).

**Proposition 3.2.12.** — [DeLi80] Let us consider a W-system \( k[[x]] \).

i) For any \( n \), \( k[[x_1, \ldots, x_n]] \) is a Noetherian Henselian regular local ring.

ii) If \( f \in k[[x_1, \ldots, x_n, y_1, \ldots, y_m]] \) and \( g := (g_1, \ldots, g_m) \in (p, x)k[[x_1, \ldots, x_n]] \), then \( f(x, g(x)) \in k[[x]] \).

iii) If \( f \in k[[x]] \), then \( \frac{\partial f}{\partial x_i} \in k[[x]] \).
iv) If $k[[x_1,\ldots,x_n]]$ is a family of rings satisfying i)-iv) of Definition 3.2.11 and if all these rings are excellent, then they satisfy v) and vi) of Definition 3.2.11.

Proof. — All these assertions are proven in Remark 1.3 [DeLi80], except iv). Thus we prove here iv): let us assume that $\text{char}(k) = p > 0$ and let $y \in (p,x)\hat{k}[x]^m$. Let us denote by $I$ the kernel of the $k[[x]]$-morphism $k[[x,y]] \to \hat{k}[x]$ defined by the substitution of $y$ for $y$ and let us assume that $I \cap k[[y]] \neq (0)$. Since $k[[x]]$ is excellent, the morphism $k[[x]] \to \hat{k}[x]$ is regular. Thus $\text{Frac}(\hat{k}[x])$ is a separable extension of $\text{Frac}(k[[x]])$, but $\text{Frac}(k[[x,y]])$ is a separable field extension. This implies that $\text{Frac}(\hat{k}[x]) \to \text{Frac}(\frac{k[[x,y]]}{I})$ is a separable field extension. This proves that Property v) of Definition 3.2.11 is satisfies. The proof that Property vi) of Definition 3.2.11 is satisfied is identical. 

Example 3.2.13. — We give here few examples of Weierstrass systems:

i) If $k$ is a field or a complete discrete valuation ring, the family $k[[x_1,\ldots,x_n]]$ is a W-system over $k$ (using Proposition 3.2.12 iv) since complete local rings are excellent rings).

ii) Let $k\langle x_1,\ldots,x_n \rangle$ be the Henselization of the localization of $k[x_1,\ldots,x_n]$ at the maximal ideal $(x_1,\ldots, x_n)$ where $k$ is a field or an excellent discrete valuation ring. Then, for $n \geq 0$, the family $k\langle x_1,\ldots,x_n \rangle$ is a W-system over $k$ (using Proposition 3.2.12 iv) since the Henselization of an excellent local ring is still excellent - see Proposition 3.8.17).

iii) The family $k\{x_1,\ldots,x_n\}$ (the ring of convergent power series in $n$ variables over a valued field $k$) is a W-system over $k$.

iv) The family of Gevrey power series in $n$ variables over a valued field $k$ is a W-system [Br86].

Then we have the following Approximation result (the case of $k\langle x \rangle$ where $k$ is a field or a discrete valuation ring is proven in [Ar69], the general case is proven in [DeLi80]):

Theorem 3.2.14. — [Ar69] [DeLi80] Let $k[[x]]$ be a W-system over $k$, where $k$ is a field or a discrete valuation ring with prime $p$. Let $f \in k[[x,y]]^r$ and $\hat{y} \in (p,x)\hat{k}[x]^m$ satisfy

$$f(x,\hat{y}) = 0.$$
Then, for any $c \in \mathbb{N}$, there exists a convergent power series solution $\tilde{y} \in (p, x)k[[x]]^m$,

$$f(x, \tilde{y}) = 0$$

such that $\tilde{y} - \hat{y} \in (p, x)^c$.

Let us mention that Theorem 3.2.8 extends also for Weierstrass systems (see [Ron10b]).

3.2.3. Néron’s desingularization and Popescu’s Theorem. — During the 70’s and the 80’s one of the main goals about Artin Approximation Problem was to find necessary and sufficient conditions on a local ring $A$ for it having the Artin Approximation Property, i.e. such that the set of solutions in $A^m$ of any system of algebraic equations $(S)$ in $m$ variables with coefficients in $A$ is dense for the Krull topology in the set of solutions of $(S)$ in $\hat{A}^m$. Let us recall that the Krull topology on $A$ is the topology induced by the following norm:

$$|a| := e^{-\text{ord}(a)}$$

for all $a \in A \setminus \{0\}$. The problem was to find a way of proving approximation results without using Weierstrass Division Theorem.

Remark 3.2.15. — Let $P(y) \in A[y]$ satisfy $P(0) \in m_A$ and $\frac{\partial P}{\partial y}(0) \notin m_A$. Then, by the Implicit Function Theorem for complete local rings, $P(y)$ has a unique root in $\hat{A}$ equal to $0$ modulo $m_A$. Thus if we want being able to approximate roots of $P(y)$ in $\hat{A}$ by roots of $P(y)$ in $A$, a necessary condition is that the root of $P(y)$ constructed by the Implicit Function Theorem is in $A$. Thus it is clear that if a local ring $A$ has the Artin Approximation Property then $A$ is necessarily Henselian.

In fact M. Artin conjectured that a sufficient condition would be that $A$ is an excellent Henselian local ring (Conjecture (1.3) [Ar70]). The idea to prove this conjecture is to generalize Płoski’s Theorem 3.2.8 and a theorem of desingularization of A. Néron [Né64]. This generalization is the following (for the definitions see Appendix 3.7):

Theorem 3.2.16. — [Po85] [Po86] Let $\varphi : A \rightarrow B$ be a regular morphism of local Noetherian rings, $C$ a finitely generated $A$-algebra and $\psi : C \rightarrow B$ a morphism of $A$-algebras. Then $\psi$ factors through a finitely generated $A$-algebra $D$ which is smooth over $A$:

$$\begin{array}{ccc}
A & \xrightarrow{\varphi} & B \\
\downarrow{\psi} & & \downarrow \\
C & \longrightarrow & D
\end{array}$$

Historically this theorem has been proven by A. Néron [Né64] when $A$ and $B$ are discrete valuation rings. Then several authors gave proofs of particular cases (see for instance [Po80], [Br83b] [ArDe83], [ArRo88], or [Rot87] - in this last paper the result is proven in the equicharacteristic zero case) until D.
Theorem 3.2.17. — Let \((A, I)\) be an excellent Henselian pair. Let \(f(y) \in A[y]^\circ\) and \(\hat{y} \in \hat{A}^m\) satisfy \(f(\hat{y}) = 0\). Then, for any \(c \in \mathbb{N}\), there exists \(\hat{y} \in A^m\) such that \(\hat{y} - \hat{y} \in I^c\hat{A}\), and \(f(\hat{y}) = 0\).

Proof. — The proof goes as follows: let us denote \(C := \frac{A[y]}{J}\) where \(J\) is the ideal generated by \(f_1, \ldots, f_r\). The formal solution \(\hat{y} \in \hat{A}\) defines an \(A\)-morphism \(\hat{\varphi} : C \rightarrow \hat{A}\) (see Remark 3.2.10). By Theorem 3.2.16, since \(A \rightarrow \hat{A}\) is regular (Example 3.7.4), there exists a smooth \(A\)-algebra \(D\) factorizing this morphism.

After some technical reductions we may assume that the morphism \(A \rightarrow D\) decomposes as \(A \rightarrow A[z] \rightarrow D\) where \(z = (z_1, \ldots, z_s)\) and \(A[z] \rightarrow D\) is standard étale. Let us choose \(\hat{z} \in \hat{A}^s\) such that \(\hat{z} - \hat{z} \in m_A^e\hat{A}^s\) (\(\hat{z}\) is the image of \(z\) in \(\hat{A}^s\)). This defines a morphism \(A[z] \rightarrow A\). Then \(A \rightarrow (\frac{D}{(z_1 - z_1, \ldots, z_s - z_s)})\) is standard étale and admits a section in \(\frac{A}{m_A^e}\). Since \(A\) is Henselian, this section lifts to a section in \(A\) by Proposition 3.8.9. This section composed with \(A[z] \rightarrow A\) defines an \(A\)-morphism \(D \rightarrow A\), and this latter morphism composed with \(C \rightarrow D\) yields a morphism \(\tilde{\varphi} : C \rightarrow A\) such that \(\tilde{\varphi}(z_i) - \tilde{\varphi}(z_i) \in m_A^e\hat{A}\) for \(1 \leq i \leq m\). 

Remark 3.2.18. — Let \((A, I)\) be a Henselian pair and let \(J\) be an ideal of \(A\). By applying this result to the Henselian pair \((\frac{J}{I}, \frac{J}{I})\) we can prove the following (using the notation of Theorem 3.2.17): if \(f(\hat{y}) \in J\hat{A}\) then there exists \(\hat{y} \in A^m\) such that \(f(\hat{y}) \in J\) and \(\hat{y} - \hat{y} \in I^cA\).

Remark 3.2.19. — In [Rot90], C. Rotthaus proves the converse of Theorem 3.2.17 in the local case: if \(A\) is a Noetherian local ring that satisfies Theorem 3.2.17, then \(A\) is excellent. In particular Weierstrass systems are excellent local rings. Previously this problem had been studied in [GiPo81] and [Br83a].

Remark 3.2.20. — Let \(A\) be a Noetherian ring and \(J\) be an ideal of \(A\). If we assume that \(f_1(y), \ldots, f_r(y) \in A[y]\) are linear, then Theorem 3.2.17 may be proven easily in this case since \(A \rightarrow \hat{A}\) is flat (see Example 3.1.3). The proof of this flatness result uses the Artin-Rees Lemma.
Example 3.2.21. — If $A$ is an excellent integral local domain let us denote by $A^h$ its Henselization. Then $A^h$ is the ring of algebraic elements of $\hat{A}$ over $A$. In particular, if $k$ is a field then $k(x)$ is the ring of formal power series which are algebraic over $k[x]$.

Indeed $A \to A^h$ is a filtered limit of algebraic extensions, thus $A^h$ is a subring of the ring of algebraic elements of $\hat{A}$ over $A$.

On the other hand if $f \in \hat{A}$ is algebraic over $A$, then $f$ satisfies an equation
\[ a_0 f^d + a_1 f^{d-1} + \cdots + a_d = 0 \]
where $a_i \in A$ for all $i$. Thus for $c$ large enough there exists $\tilde{f} \in A^h$ such that $\tilde{f}$ satisfies the same polynomial equation and $\tilde{f} - f \in m_A^c$ (by Theorem 3.2.17 and Theorem 3.8.17). Since $\bigcap m_A^c = (0)$ and a polynomial equation has a finite number of roots, this proves that $\tilde{f} = f$ for $c$ large enough and $f \in A^h$.

Example 3.2.22. — The strength of this result comes from the fact that it applies to rings that do not satisfy the Weierstrass Preparation Theorem. For example Theorem 3.2.17 applies to the local ring $B = A[x_1, \ldots, x_n]$ where $A$ is an excellent Henselian local ring (the main example is $A = k[[t]](x)$ where $t$ and $x$ are multivariables). Indeed, this ring is the Henselization of $A[x_1, \ldots, x_n][x_1 + (x_1, \ldots, x_n)]$. Thus $B$ is an excellent local ring by Example 3.7.4 and Proposition 3.8.17.

This case was the main motivation of D. Popescu for proving Theorem 3.2.16 (see also [Ar70]), since this case implies a nested Artin Approximation result (see Theorem 3.5.8).

Previous particular cases of this application had been studied before: see [PfPo81] for a direct proof that $V[x_1, x_2]$ satisfies Theorem 3.2.17 when $V$ is a complete discrete valuation ring, and [BDL83] for the ring $k[x_1, x_2][x_3, x_4, x_5]$.

Hint of the proof of Theorem 3.2.16 — Let $A$ be a Noetherian ring and $C$ be a $A$-algebra of finite type, $C = A[y_1, \ldots, y_m]$ with $I = (f_1, \ldots, f_r)$. We denote by $\Delta_g$ the ideal of $A[y]$ generated by the $h \times h$ minors of the Jacobian matrix
\[ \left( \frac{\partial g_j}{\partial y_i} \right)_{1 \leq i \leq h, 1 \leq j \leq m} \] for $g := (g_1, \ldots, g_h) \subset I$. We define the ideal
\[ H_{C/A} := \sqrt{\sum_g \Delta_g((g) : I)C} \]
where the sum runs over all $g := (g_1, \ldots, g_h) \subset I$ and $h \in \mathbb{N}$. This ideal is independent of the presentation of $C$ and it defines the singular locus of $C$ over $A$.

Lemma 3.2.23. — For any $p \in \text{Spec}(C)$, $C_p$ is smooth over $A$ if and only if $H_{C/A} \not\subset p$. 

We have the following property:

**Lemma 3.2.24.** — Let $C$ and $C'$ be two $A$-algebras of finite type and let $A \rightarrow C \rightarrow C'$ be two morphisms of $A$-algebras. Then $H_{C'/C} \cap \sqrt{H_{C'/A} C'} = H_{C'/C} \cap \sqrt{H_{C'/A}}$.

The idea of the proof of Theorem 3.2.16 is the following: if $H_{C/A} B \neq B$, then we replace $C$ by a $A$-algebra of finite type $C'$ such that $H_{C/A} B$ is a proper sub-ideal of $H_{C'/A} B$. Using the Noetherian assumption, after a finite number we have $H_{C/A} B = B$. Then we use the following proposition:

**Proposition 3.2.25.** — Using the notation of Theorem 3.2.16, let us assume that $H_{C/A} B = B$. Then $\psi$ factors as in Theorem 3.2.16.

**Proof of Proposition 3.2.25.** — Let $(c_1, ..., c_s)$ be a system of generators of $H_{C/A}$. Then $1 = \sum_{i=1}^{s} b_i \psi(c_i)$ for some $b_i$’s in $B$. Let us define

$$D := \frac{C[z_1, ..., z_s]}{(1 - \sum_{i=1}^{s} c_i z_i)}.$$

We construct a morphism of $C$-algebra $D \rightarrow B$ by sending $z_i$ onto $b_i$, $1 \leq i \leq s$. It is easy to check $D_{c_i}$ is a smooth $C$-algebras, thus $c_i \in H_{D/C}$ by Lemma 3.2.23, and $H_{C/A} D \subset H_{D/C}$. By Lemma 3.2.24, since $1 \in H_{C/A} D$, we see that $1 \in H_{D/A}$. By Lemma 3.2.23 this proves that $D$ is a smooth $A$-algebra.

Now to increase the size of $H_{C/A} B$ we use the following proposition:

**Proposition 3.2.26.** — Using the notation of Theorem 3.2.16, let $p$ be a minimal prime ideal of $H_{C/A} B$. Then there exist a factorization of $\psi : C \rightarrow D \rightarrow B$ such that $D$ is finitely generated over $A$ and $\sqrt{H_{C/A} B} \subsetneq \sqrt{H_{D/A} B} \subsetneq p$.

The proof of Proposition 3.2.26 is done by induction on $\text{height}(p)$. Thus there are two things to prove: first the case $\text{ht}(p) = 0$ which is equivalent to prove Theorem 3.2.16 for Artinian rings, then the reduction $\text{ht}(p) = k+1$ to the case $\text{ht}(p) = k$. This last case is quite technical, even in the equicharacteristic zero case (i.e. when $A$ contains $\mathbb{Q}$, see [Qu97] for a good presentation of this case). In the case $A$ does not contain $\mathbb{Q}$ there appear more problems due to the existence of inseparable extensions of residue fields. In this case the André homology is the good tool to handle these problems (see [Sw98]).
3.3. Strong Artin Approximation

We review here results about the Strong Approximation Property. There is clearly two different cases: the first case is when the base ring is a discrete valuation ring (where life is easy!) and the second case is the general case (where life is less easy).

3.3.1. Greenberg’s Theorem: the case of a discrete valuation ring.

— Let $V$ be a Henselian discrete valuation ring, $\mathfrak{m}_V$ its maximal ideal and $\mathbb{K}$ be its field of fractions. Let us denote by $\hat{V}$ the $\mathfrak{m}_V$-adic completion of $V$ and by $\hat{\mathbb{K}}$ its field of fractions. If $\text{char}(\mathbb{K}) > 0$, let us assume that $\mathbb{K} \rightarrow \hat{\mathbb{K}}$ is a separable field extension (in this case this is equivalent to $V$ being excellent, see Example 3.7.2 iii) and Example 3.7.4 iv)).

**Theorem 3.3.1 (Greenberg’s Theorem).** — [Gre66] If $f(y) \in V[y]^r$, then there exist $a, b \geq 0$ such that

$$\forall c \in \mathbb{N} \forall \overline{y} \in V^m \text{ such that } f(\overline{y}) \in \mathfrak{m}_V^{ac+b},$$

$$\exists \tilde{y} \in V^m \text{ such that } f(\tilde{y}) = 0 \text{ and } \tilde{y} - y \in \mathfrak{m}_V^c.$$

**Sketch of proof.** — We will give the proof in the case $\text{char}(\mathbb{K}) = 0$. The result is proven by induction on the height of the ideal generated by $f_1(y), ..., f_r(y)$. Let us denote by $I$ this ideal. We will denote by $\nu$, the $\mathfrak{m}_V$-adic order on $V$ which is a valuation by assumption.

There exists an integer $e \geq 1$ such that $\sqrt{I^e} \subset I$. Then $f(\overline{y}) \in \mathfrak{m}_V^\nu$ for all $f \in I$ implies that $f(\overline{y}) \in \mathfrak{m}_V^\nu$ for all $f \in \sqrt{I}$ since $V$ is a valuation ring. Moreover if $\sqrt{I} = \mathcal{P}_1 \cap \cdots \cap \mathcal{P}_s$ is prime decomposition of $\sqrt{I}$, then $f(\overline{y}) \in \mathfrak{m}_V^\nu$ for all $f \in \sqrt{I}$ implies that $f(\overline{y}) \in \mathfrak{m}_V^\nu$ for all $f \in \mathcal{P}_i$ for some $i$. This allows us to assume that $I$ is a prime ideal of $V[y]$.

Let $h$ be the height of $I$. If $h = m + 1$, then $I$ is a maximal ideal of $V[y]$ and thus it contains some non zero element of $V$ denoted by $v$. Then there does not exist $\overline{y} \in V^m$ such that $f(\overline{y}) \in \mathfrak{m}_V^{\nu(v)+1}$ for all $f \in I$. Thus the theorem is true for $a = 0$ and $b = \nu(v) + 1$.

Let us assume that the theorem is proven for ideals of height $h + 1$ and let $I$ be a prime ideal of height $h$. As in the proof of Theorem 3.2.1 we may assume that $r = h$ and that the determinant of the Jacobian matrix of $f$, denoted by $\delta$, is not in $I$. Let us denote $J := I + (\delta)$. Since $\text{ht}(J) = h + 1$, by the inductive hypothesis, there exist $a, b \geq 0$ such that

$$\forall c \in \mathbb{N} \forall \overline{y} \in V^m \text{ such that } f(\overline{y}) \in \mathfrak{m}_V^{ac+b} \forall f \in J$$

$$\exists \tilde{y} \in V^m \text{ such that } f(\tilde{y}) = 0 \forall f \in J \text{ and } \tilde{y}_j - y_j \in \mathfrak{m}_V, \ 1 \leq j \leq m.$$
Then let \( c \in \mathbb{N} \) and \( \bar{y} \in V^m \) satisfy \( f(\bar{y}) \in m^{(2a+1)c+2b}_V \) for all \( f \in I \). If \( \delta(\bar{y}) \in m^{ac+b}_V \), then \( f(\bar{y}) \in m^{ac+b}_V \) for all \( f \in J \) and the result is proven by the inductive hypothesis. If \( \delta(\bar{y}) \notin m^{ac+b}_V \), then \( f_i(\bar{y}) \in (\delta(\bar{y}))^2 m^{c}_V \) for \( 1 \leq i \leq r \). Then the result comes from the following result.

**Theorem 3.3.2 (Tougeron’s Implicit Function Theorem)**

Let \( A \) be a Henselian local ring and \( f(y) \in A[y]^r, y = (y_1, ..., y_m), m \geq r \). Let \( \delta(x,y) \) be a \( r \times r \) minor of the Jacobian matrix \( \frac{\partial (f_1, ..., f_r)}{\partial (y_1, ..., y_m)} \). Let us assume that there exists \( y \in A^m \) such that \( f_i(y) \in (\delta(y))^2 m^{c}_A \) for all \( 1 \leq i \leq r \) and for some \( c \in \mathbb{N} \). Then there exists \( \tilde{y} \in A^m \) such that \( f_i(\tilde{y}) = 0 \) for all \( 1 \leq i \leq r \), and \( \tilde{y} - y \in (\delta(y))m^{c}_A \).

**Proof.** — The proof is completely similar to the proof of Theorem 3.2.5.

In fact we can prove the following result whose proof is identical to the proof of Theorem 3.3.1:

**Theorem 3.3.3.** — [Sc83] Let \( V \) be a complete discrete valuation ring and \( f(y,z) \in V[y][z]^r \), where \( z := (z_1, ..., z_s) \). Then there exist \( a, b \geq 0 \) such that

\[
\forall c \in \mathbb{N} \ \forall \bar{y} \in (m_V V)^m, \ \forall \bar{z} \in V^s \text{ such that } f(\bar{y}, \bar{z}) \in m^{ac+b}_V
\]

\[
\exists \tilde{y} \in (m_V V)^m, \ \exists \tilde{z} \in V^s \text{ such that } f(\tilde{y}, \tilde{z}) = 0 \text{ and } \tilde{y} - y, \tilde{z} - z \in m^{c}_V.
\]

**Remark 3.3.4.** — M. Greenberg proved this result in order to study \( C_i \) fields. Previous results about \( C_i \) fields had been already been studied, in particular by S. Lang in [Lan52] where appeared for the first time a particular case of Artin Approximation Theorem (see Theorem 11 and its corollary in [Lan52]).

**Remark 3.3.5.** — In the case \( f(y) \) has no solution in \( V \), we can choose \( a = 0 \) and Theorem 3.3.1 asserts there exists a constant \( b \) such that \( f(y) \) has no solution in \( V^m \).

**Remark 3.3.6.** — The valuation \( \nu \) of \( V \) defines a ultrametric norm on \( K \): we define it as

\[
\left| \frac{y}{z} \right| := e^{\nu(z) - \nu(y)}, \ \forall y, z \in V \setminus \{0\}.
\]

This norm defines a distance on \( V^m \), for any \( m \in \mathbb{N}^* \), denoted by \( d(.,.) \) and defined by

\[
d(y, z) := \max_{k=1}^m |y_k - z_k|.
\]
Then Theorem 3.3.1 can be reformulated as a Łojasiewicz Inequality (see Te12):  

\[ \exists a \geq 1, \ C > 0 \text{ s.t. } |f(\overline{y})| \geq Cd(f^{-1}(0), \overline{y})^a \forall \overline{y} \in V^m. \]

This Łojasiewicz Inequality is well known for algebraic or analytic functions and Theorem 3.3.1 can be seen as a generalization of this Łojasiewicz Inequality for algebraic or analytic functions defined over \( V \). If \( V = k[t] \) where \( k \) is a field, there is very few results known about the geometry of algebraic varieties defined over \( V \). It is a general problem to extend classical results of differential or analytic geometry over \( \mathbb{R} \) or \( \mathbb{C} \) to this setting. See for instance [HaMü94], [BH10] (extension of Rank Theorem), [Reg06] or [FBPP12b] (Extension of Curve Selection Lemma), [Hic05] for some results in this direction.

For any \( c \in \mathbb{N} \) let us denote by \( \beta(c) \) the smallest integer such that:

for all \( \overline{y} \in V^m \) such that \( f(\overline{y}) \in (x)^{\beta(c)} \), there exists \( \tilde{y} \in V^m \) such that \( f(\tilde{y}) = 0 \) and \( \tilde{y} - y \in (x)^c \). Greenberg’s Theorem asserts that such a function \( \beta : \mathbb{N} \rightarrow \mathbb{N} \) exists and that it is bounded by an affine function. We call this function \( \beta \) the Greenberg function of \( f \). We can remark that the Greenberg function is an invariant of the integral closure of the ideal generated by \( f_1, ..., f_r \):

**Lemma 3.3.7.** — Let us consider \( f(y) \in V[y]^r \) and \( g(y) \in V[y]^q \). Let us denote by \( \beta_f \) and \( \beta_g \) their Greenberg functions. Let \( I \) (resp. \( J \)) be the ideal of \( V[y] \) generated by \( f_1(y), ..., f_r(y) \) (resp. \( g_1(y), ..., g_q(y) \)). If \( I = J \) then \( \beta_f = \beta_g \). The same is true for Theorem 3.3.3.

**Proof.** — Let \( \mathcal{I} \) be an ideal of \( V \) and \( \overline{y} \in V^m \). We remark that

\[ f_1(\overline{y}), ..., f_r(\overline{y}) \in \mathcal{I} \iff g(\overline{y}) \in \mathcal{I} \forall g \in I. \]

Then by replacing \( \mathcal{I} \) by \( (0) \) and \( m^c \), for all \( c \in \mathbb{N} \), we see that \( \beta_f \) depends only on \( I \).

Now, for any \( c \in \mathbb{N} \), we have:

\[ g(\overline{y}) \in m^c_I \forall g \in I \iff \nu(g(\overline{y})) \geq c \forall g \in I \]

\[ \iff \nu(g(\overline{y})) \geq c \forall g \in \mathcal{I} \]

\[ \iff g(\overline{y}) \in m^c_{\mathcal{I}} \forall g \in \mathcal{I}. \]

Thus \( \beta_f \) depends only on \( \mathcal{I} \).

In general, it is a difficult problem to compute the Greenberg function of an ideal \( I \). It is even a difficult problem to bound this function in general. We can give few results giving some informations about the Greenberg functions. First
of all let us remark that, in the proof of Theorem 3.3.1, we proved a particular case of the following inequality:

$$\beta_I(c) \leq 2\beta_J(c) + c, \quad \forall c \in \mathbb{N}$$

where $J$ is the Jacobian Ideal of $I$ (for a precise definition of the Jacobian Ideal in general and a general proof of this inequality let see [Elk73]). The coefficient 2 comes from the use of Tougeron’s Implicit Function Theorem. We can sharpen this bound in the following particular case:

**Theorem 3.3.8.** — [Hic93] Let $k$ be an algebraically closed field of characteristic zero and $V := k[[t]]$ where $t$ is a single variable. Let $f(y) \in V[y]$ be one power series. Let us denote by $J$ the ideal of $V[y]$ generated by $f(y), \frac{\partial f}{\partial y}(y), \frac{\partial^2 f}{\partial y^2}(y), \ldots, \frac{\partial^m f}{\partial y^m}(y)$, and let us denote by $\beta_f$ the Greenberg function of $(f)$ and by $\beta_J$ the Greenberg function of $J$. Then

$$\beta_f(c) \leq \beta_J(c) + c, \quad \forall c \in \mathbb{N}.$$ 

This bound may be used to find sharp bounds of some Greenberg functions (see Remark 3.3.10).

On the other hand we can describe the behaviour of $\beta$ in the following case:

**Theorem 3.3.9.** — [De84][DeLo99] Let $V$ be $\mathbb{Z}_p$ or a Henselian discrete valuation ring whose residue field is an algebraically closed field of characteristic zero. Let us denote by $\mathfrak{m}_V$ the maximal ideal of $V$. Let us denote by $\beta$ the Artin function of $f(y) \in V[y]^r$. Then there exists a finite partition of $\mathbb{N}$ in congruence classes such that on each such congruence class the function $c \mapsto -\beta(c)$ is linear for $c$ large enough.

**Hints on the proof in the case the residue field has characteristic zero**

Let us consider the following first order language of three sorts:

1) the field $(K := \text{Frac}(V), +, \times, 0, 1)$
2) the group $(\mathbb{Z}, +, <, \equiv_d (\forall d \in \mathbb{N}^*), 0) (\equiv_d$ is the relation $a \equiv_d b$ if and only if $a - b$ is divisible by $d$ for $a, b \in \mathbb{Z})$
3) the residue field $(k := \text{Frac}(\frac{V}{\mathfrak{m}_V}), +, \times, 0, 1)$

with both following functions:

a) $\nu : K \rightarrow \mathbb{Z}^*$

b) $ac : K \rightarrow k$ ("angular component")

The function $\nu$ is the valuation of the valuation ring $V$. The function $ac$ may be characterized by axioms, but here let us just give an example: let us assume that $V = k[[t]]$. Then $ac$ is defined by $ac(0) = 0$ and $ac(\sum_{n=n_0}^\infty a_nt^n) = a_{n_0}$ if $a_{n_0} \neq 0$. 


3.3. STRONG ARTIN APPROXIMATION

The second sort \((\mathbb{Z}, +, <, \equiv, 0)\) admits elimination of quantifiers \((\text{Pr29})\) and the elimination of quantifiers of \((k, +, \times, 0, 1)\) is a classical result of Chevalley. J. Pas proved that the first sort language and the three sorted language admits elimination of quantifiers \([\text{Pas89}]\). This means that any subset of \(K^{n_1} \times \mathbb{Z}^{n_2} \times k^{n_3}\) defined by a first order formula in this three sorts language (i.e. a logical formula involving 0, 1, +, \times (but not \(a \times b\) where \(a\) and \(b\) are integers), (, ), =, <, \&, \lor, \neg, \forall, \exists, \nu, ac\), and variables for elements of \(K, \mathbb{Z}\) and \(k\) may be defined by a formula involving the same symbols except \(\forall, \exists\).

Then we see that \(\beta\) is defined by the following logical sentence:

\[
\forall c \in \mathbb{N} \forall \bar{y} \in K^m (\nu(f(\bar{y})) \geq \beta(c)) \land (\nu(\bar{y}) \geq 0) \exists \tilde{y} \in K^m (f(\tilde{y}) = 0 \land \nu(\tilde{y} - y) \geq c)] \\
\land [\forall c \in \mathbb{N} \exists \bar{y} \in K^m (\nu(f(\bar{y})) \geq \beta(c) + 1) \land (\nu(\bar{y}) \geq 0) \\
\neg \exists \tilde{y} \in K^m (f(\tilde{y}) = 0 \land \nu(\tilde{y} - y) \geq c)]
\]

Applying the latter elimination of quantifiers result we see that \(\beta(c)\) may be defined without \(\forall\) and \(\exists\). Thus \(\beta(c)\) is defined by a formula using \(+, <, \equiv\) (for a finite set of integers \(d\)). This proves the result.

The case where \(V = \mathbb{Z}_p\) needs more work since the residue field of \(\mathbb{Z}_p\) is not algebraically closed, but the idea is the same.

\[\text{Remark 3.3.10.} \quad \text{When } V = \mathbb{C}\{t\}, \text{ } t \text{ being a single variable, it is tempting to link together the Greenberg function of a system of equations with coefficients in } V \text{ and some geometric invariants of the germ of complex set defined by this system of equations. This has been done in several cases:}
\]

i) In \([\text{El89}]\), a bound (involving the multiplicity and the Milnor number) of the Greenberg function is given when the system of equations defines a curve in \(\mathbb{C}^m\).

ii) Using Theorem 3.3.8 \([\text{Hic93}]\) gives the following bound of the Greenberg function \(\beta\) of a germ of complex hypersurface with an isolated singularity: 

\(\beta(c) \leq |\lambda c| + c\) for all \(c \in \mathbb{N}\), and this bound is sharp for plane curves. Here \(\lambda\) denotes the Łojasiewicz exponent of the germ, i.e.

\[\lambda := \inf \{\theta \in \mathbb{R} \mid \exists C > 0 \exists U \text{ neighborhood of } 0 \text{ in } \mathbb{C}^m, \]

\[|f(z)| + \left| \frac{\partial f}{\partial z_1}(z) \right| + \cdots + \left| \frac{\partial f}{\partial z_m}(z) \right| \geq C|z|^\theta \forall z \in U\} .\]

iii) \([\text{Hic04}]\) makes the complete computation of the Greenberg function of a branch of plane curve and proves that it is a topological invariant. This computation has been done for several branches in \([\text{Sa10}]\). Some particular cases depending on the Newton polygon of the plane curve singularity are computed in \([\text{Wa78}]\).
Finally we mention the following recent result that extends Theorem 3.3.1 to non-Noetherian valuation rings and whose proof is based on ultraproducts methods used in [BDLvdD79] to prove Theorem 3.3.1 (see 3.3.3):

**Theorem 3.3.11.** — [M-B11] Let $V$ be a Henselian valuation ring and $\nu : V \to \Gamma$ its associated valuation. Let us denote by $\hat{V}$ its $m_V$-adic completion, $\mathbb{K} := \text{Frac}(V)$ and $\hat{\mathbb{K}} := \text{Frac}(\hat{V})$. Let us assume that $\mathbb{K} \to \hat{\mathbb{K}}$ is a separable extension. Then for any $f(y) \in V[y]^r$ there exist $a \in \mathbb{N}$, $b \in \Gamma^+$ such that

$$\forall c \in \Gamma \forall y \in V^m (\nu(f(y)) \geq ac + b) \implies \exists \tilde{y} \in V^m (f(\tilde{y}) = 0 \land \nu(\tilde{y} - y) \geq c).$$

**3.3.2. Strong Artin Approximation Theorem: the general case.** — In the general case (when $V$ is not a valuation ring), there still exists an approximation function $\beta$. We have the following results:

**Theorem 3.3.12.** — [Ar68, BDLvdD79] Let $k$ be a field. For all $n, m, d \in \mathbb{N}$, there exists a function $\beta_{n,m,d} : \mathbb{N} \to \mathbb{N}$ such that the following holds:

Set $x := (x_1, ..., x_n)$ and $y := (y_1, ..., y_m)$. Then for all $f(x, y) \in k[x, y]^r$ of total degree $\leq d$, for all $c \in \mathbb{N}$, for all $\tilde{y}(x) \in k[[x]]^m$ such that

$$f(x, \tilde{y}(x)) \in (x)^{\beta_{n,m,d}(c)},$$

there exists $\tilde{y}(x) \in k[[x]]^m$ such that $f(\tilde{y}(x)) = 0$ and $\tilde{y}(x) - y(x) \in (x)^c$.

**Remark 3.3.13.** — By following the proof of M. Artin, D. Lascar proved that there exists a recursive function $\beta$ that satisfies the conclusion of Theorem 3.3.12 [Las78]. But the proof of Theorem 3.3.12 uses a double induction on the height of the ideal (like in Theorem 3.3.1) and on $n$ (like in Theorem 3.2.1). In particular, in order to apply the Jacobian Criterion, we need to work with prime ideals, and replace the original ideal $I$ generated by $f_1, ..., f_r$ by one of its associated prime and then make a reduction to $n - 1$ variables. But the bounds of the degree of the generators of such associated prime may be very large compared to the degree of the generators of $I$. This is essentially the reason why the proof of this theorem does not give much more information about the growth of $\beta$ than Lascar’s result.

**Theorem 3.3.14.** — [PiPo75, Po86] Let $A$ be a complete local ring whose maximal ideal is denoted by $m_A$. Let $f(y, z) \in A[[y]][[z]]^r$, with $z := (z_1, ..., z_s)$. Then there exists a function $\beta : \mathbb{N} \to \mathbb{N}$ such that the following holds:

For any $c \in \mathbb{N}$ and any $\tilde{y} \in (m_A, A)^m$ and $\tilde{z} \in A^s$ such that $f(\tilde{y}, \tilde{z}) \in m_A^{\beta(c)}$, there exists $\tilde{y} \in (m_A, A)^m$ and $\tilde{z} \in A^s$ such that $f(\tilde{y}, \tilde{z}) = 0$ and $\tilde{y} - \tilde{y}, \tilde{z} - \tilde{z} \in m_A^c$. 
Example 3.3.15. — [Sp94] Set $f(x_1, x_2, y_1, y_2) := x_1y_1^2 - (x_1 + x_2)y_2^2$. Set
\[
\sqrt{1 + t} = 1 + \sum_{n \geq 1} a_n t^n \in \mathbb{C}[t]
\]
be the power series such that $\sqrt{1 + t^2} = 1 + t$. For any $c \in \mathbb{N}$ set $y_2^{(c)}(x) := x_1^c$ and $y_1^{(c)}(x) := x_1^c + \sum_{n=1}^{c} a_n x_1^{c-n} x_2^n$. Then
\[
f(x_1, x_2, y_1^{(c)}(x), y_2^{(c)}(x)) \in (x_2)^c.
\]
On the other side the equation $f(x_1, x_2, y_1(x), y_2(x)) = 0$ has no other solution $(y_1(x), y_2(x)) \in \mathbb{k}[x]^2$ but $(0, 0)$. This proves that Theorem 3.3.16 is not valid for general Henselian pairs since $(\mathbb{k}[x_1, x_2], (x_2))$ is a Henselian pair.

Let us notice that L. Moret-Bailley proved that if a pair $(A, I)$ satisfies Theorem 3.3.16 then $A$ has to be an excellent Henselian local ring [M-B07]. It is still open to know under which conditions on $I$ the pair $(A, I)$ satisfies Theorem 3.3.16 when $A$ is an excellent Henselian local ring.

Corollary 3.3.16. — [PiPo75] [Po86] Let $A$ be an excellent Henselian local ring whose maximal ideal is denoted by $m_A$ and let $f(y) \in A[y]^r$. Then there exists a function $\beta : \mathbb{N} \to \mathbb{N}$ such that:
\[
\forall c \in \mathbb{N}, \forall \overline{y} \in A^m \text{ such that } f(\overline{y}) \in m_A^{\beta(c)}
\]
\[
\exists \overline{y} \in A^m \text{ such that } f(\overline{y}) = 0 \text{ and } \overline{y} - \overline{y} \in m_A^{\beta(c)}.
\]

Corollary 3.3.17. — [Wa75] [DeLi80] Let $\mathbb{k}[\lbrack x \rbrack]$ be a W-system over $\mathbb{k}$, where $\mathbb{k}$ is a field or a discrete valuation ring with prime $\mathfrak{p}$. Let $f(x, y) \in \mathbb{k}[\lbrack x, y \rbrack]^r$. Then there exists a function $\beta : \mathbb{N} \to \mathbb{N}$ such that for any $c \in \mathbb{N}$ and any $\overline{y} \in (\mathfrak{p}, x)\mathbb{k}[\lbrack x \rbrack]^m$ such that $f(x, \overline{y}) \in (\mathfrak{p}, x)^{\beta(c)}$, there exists $\tilde{y} \in (\mathfrak{p}, x)\mathbb{k}[\lbrack x \rbrack]^m$ such that $f(x, \tilde{y}) = 0$ and $\tilde{y} - \overline{y} \in (\mathfrak{p}, x)^{\beta(c)}$.

Proof. — We first apply Theorem 3.3.14 then we apply Theorem 3.2.14. \qed

Remark 3.3.18. — As for Theorem 3.3.1 Corollary 3.3.16 implies that, if $f(y)$ has no solution in $A$, there exists a constant $c$ such that $f(y)$ has no solution in $A_{m_A}^{1/c}$.

Definition 3.3.19. — Let $f$ be as in Theorem 3.3.14 or Corollary 3.3.16. The least function $\beta$ that satisfies these theorems is called the Artin function of $f$.

Remark 3.3.20. — As before, the Artin function of $f$ depends only on the integral closure of the ideal $I$ generated by $f_1, \ldots, f_r$ (see Lemma 3.3.7).
**Remark 3.3.21.** — Let \( f(y) \in A[y]^r \) and \( y \in (m_A)^m \) satisfy \( f(y) \in m_A^r \) and let us assume that \( A \rightarrow B := \frac{A[y]}{(f(y))} \) is a smooth morphism. This morphism is local thus it splits as \( A \rightarrow C := A[z]_{m_A+(z)} \rightarrow B \) such that \( C \rightarrow B \) is étale (see Definition 3.8.5) and \( z := (z_1, ..., z_s) \). We remark that \( y \) defines a morphism of \( A \)-algebras \( \varphi : B \rightarrow A \) mod \( m_A \). Let us choose any \( \tilde{z} \in A^s \) such that \( z_i - \tilde{z}_i \in m_A^s \) for all \( 1 \leq i \leq s \) (\( z_i \) denotes the image of \( z_i \) in \( \frac{A}{m_A} \)). Then \( A \rightarrow \frac{B}{(z_1 - \tilde{z}_1, ..., z_s - \tilde{z}_s)} \) is étale and admits a section in \( \frac{A}{m_A} \). By Proposition 3.8.9 this section lifts to a section in \( A \). Thus we have a section \( B \rightarrow A \) equal to \( \varphi \) modulo \( m_A \). This proves that \( \beta(c) = c \) when \( A \rightarrow \frac{A[y]}{(f(y))} \) is smooth.

### 3.3.3. Ultraproducts and proofs of Strong Approximation type results.

Historically, M. Artin proved Theorem 3.3.12 in [Ar69] by slightly modifying the proof of Theorem 3.2.1, i.e. by an induction on \( n \) using the Weierstrass Division Theorem. Then some people tried to prove this kind of result in the same way, but this was not always easy, in particular when the base field was not a characteristic zero field (for example there is a gap in the inseparable case of [PfPo75]). Then four people introduced the use of ultraproducts to give easy proofs of this kind of Strong Approximation type results ([BDLvdD79] and [DeLi80]; see also [Po79] for the general case). The general principle is the following: ultraproducts reduce Strong Artin Approximation Problems to Artin Approximation Problems. We will present here the main ideas.

Let us start with some terminology. A **filter** \( D \) (over \( \mathbb{N} \)) is a non empty subset of \( \mathcal{P}(\mathbb{N}) \), the set of subsets of \( \mathbb{N} \), that satisfies the following properties:

\[
\begin{align*}
& a) \emptyset \notin D, \quad b) \mathcal{E}, \mathcal{F} \in D \implies \mathcal{E} \cap \mathcal{F} \in D, \quad c) \mathcal{E} \in D, \mathcal{F} \subseteq \mathcal{E} \implies \mathcal{F} \in D.
\end{align*}
\]

A filter \( D \) is **principal** if \( D = \{ \mathcal{F} \mid \mathcal{E} \subseteq \mathcal{F} \} \) for some subset \( \mathcal{E} \) of \( \mathbb{N} \). A **ultrafilter** is a filter which is maximal for the inclusion. It is easy to check that a filter \( D \) is a ultrafilter if and only if for any subset \( \mathcal{E} \) of \( \mathbb{N} \), \( D \) contains \( E \) or its complement \( \mathbb{N} - \mathcal{E} \). In the same way an ultrafilter is non-principal if and only if it contains the filter \( E := \{ \mathcal{E} \subseteq \mathbb{N} \mid \mathbb{N} - \mathcal{E} \text{ is finite} \} \). Zorn’s Lemma yields the existence of non-principal ultrafilters.

Let \( A \) be a Noetherian ring. Let \( D \) be a non-principal ultrafilter. We define the ultrapower (or ultraproduct) of \( A \) as follows:

\[
A^* := \frac{\{(a_i)_{i \in \mathbb{N}} \in \prod_i A \}}{(a_i) \sim (b_i) \iff \{i/a_i = b_i \in D\}}.
\]
We have a morphism $A \rightarrow A^*$ that sends $a$ onto the class of $(a)_i \in \mathbb{N}$. We have the following fundamental result:

**Theorem 3.3.22.** — [ChKe73] Let $L$ be a first order language, let $A$ be a structure for $L$ and let $D$ be an ultrafilter over $\mathbb{N}$. Then for any $(a_i)_{i \in \mathbb{N}} \in A^*$ and for any first order formula $\varphi(x), \varphi((a_i))$ is true in $A^*$ if and only if \( \{i \in \mathbb{N} / \varphi(a_i) \text{ is true in } A\} \in D \).

In particular we can deduce the following properties:

The ultrapower $A^*$ is equipped with a structure of commutative ring. If $A$ is a field then $A^*$ is a field. If $A$ is an algebraically closed field then $A^*$ is an algebraically closed field. If $A^*$ is a local ring with maximal ideal $m_A$ then $A^*$ is a local ring with maximal ideal $m_A^*$ defined by $(a_i), \in m_A$ if and only if \( \{i / a_i \in m_A\} \in D \). If $A$ is a local Henselian ring, then $A^*$ is a local Henselian ring. In fact all these facts are elementary and can be checked directly by hand. Elementary proofs of these results can be found in [BdLvdD79]. Nevertheless if $A$ is Noetherian, then $A^*$ is not Noetherian in general, since Noetherianity is a condition on ideals of $A$ and not on elements of $A$. For example, if $A$ is a Noetherian local ring, then \( m_{\infty}^* := \bigcap_{n \geq 0} m_A^n \neq (0) \) in general. But we have the following lemma:

**Lemma 3.3.23.** — [Po00] Let $(A, m_A)$ be a Noetherian complete local ring. Let us denote $A_1 := A/m_A^\infty$. Then $A_1$ is a Noetherian complete local ring of same dimension as $A$ and the composition $A \rightarrow A^* \rightarrow A_1$ is flat.

In fact, since $A$ is excellent and $m_A A_1$ is the maximal ideal of $A_1$, it is not difficult to prove that $A \rightarrow A_1$ is regular. Details can be found in [Po00].

Let us sketch the idea in the case of Theorem 3.3.16:

**Sketch of the proof of Theorem 3.3.16.** — Let us assume that some system of algebraic equations over an excellent Henselian local ring $A$, denoted by $f = 0$, does not satisfy Theorem 3.3.16. Using Theorem 3.2.17 we may assume that $A$ is complete. Thus it means that there exist an integer $c_0 \in \mathbb{N}$ and $\overline{y}^{(c)} \in A^m$, $\forall c \in \mathbb{N}$, such that $f(\overline{y}^{(c)}) \in m_A^c$ and there does not exists $\tilde{y}^{(c)} \in A^m$ such that $f(\tilde{y}^{(c)}) = 0$ and $\tilde{y}^{(c)} - \overline{y}^{(c)} \in m_A^c$.

Let us denote by $\overline{y}$ the image of $(\overline{y}^{(c)})_c$ in $(A^*)^m$. Since $f(y) \in A[y]^r$, we may assume that $f(y) \in A^*[y]^r$ using the morphism $A \rightarrow A^*$. Then $f(\overline{y}) \in m_{\infty}^*$. Thus $f(\overline{y}) = 0$ in $A_1$. Let us choose $c > c_0$. Since $A \rightarrow A_1$ is regular and $A$ is Henselian, following the proof of Theorem 3.2.17 for any $c \in \mathbb{N}$ there exists $\tilde{y} \in A^m$ such that $f(\tilde{y}) = 0$ and $\tilde{y} - \overline{y} \in m_A A_1$. Thus $\tilde{y} - \overline{y} \in m_A A^*$. Hence the set $\{i \in \mathbb{N} / \tilde{y} - \overline{y}^{(c)} \in m_A A^*\} \in D$ is non-empty. This is a contradiction. \( \square \)
We can also prove easily the following proposition with the help of ultraproducts:

**Proposition 3.3.24.** — Let $f(x, y) \in \mathbb{C}[x, y]^r$. For any $1 \leq i \leq m$ let $J_i$ be a subset of $\{1, \ldots, n\}$.

Let us assume that, for any $c \in \mathbb{N}$, there exist $y_i^{(c)}(x) \in \mathbb{C}[x_j, j \in J_i]$, $1 \leq i \leq m$, such that

$$f(x, y_i^{(c)}(x)) \in (x)^c.$$ 

Then there exist $\bar{y}_i(x) \in \mathbb{C}[x_j, j \in J_i]$, $1 \leq i \leq m$, such that $f(x, \bar{y}_i(x)) = 0$.

**Proof.** — Let us denote by $\bar{y} \in \mathbb{C}[x]^*$ the image of $(y_i^{(c)})_c$. Then $f(x, \bar{y}) = 0$ modulo $(x)^\infty$. It is not very difficult to check that $\mathbb{C}[x]^* / (x)^\infty \simeq \mathbb{C}^*[x]$ as $\mathbb{C}^*$-algebras. Moreover $\mathbb{C}^*$ is an $\mathbb{k}$-algebra (where $\mathbb{k}$ is the subfield of $\mathbb{C}$ generated by the coefficients of $f$), since they are algebraically closed field of same transcendence degree over $\mathbb{Q}$ and same characteristic. Then the image of $\bar{y}$ by the isomorphism yields the desired solution in $\mathbb{C}[x]$.

Let us remark that the proof of this result remains valid if we replace $\mathbb{C}$ by any algebraically closed field $\mathbb{k}$ whose cardinal is strictly greater that the cardinal of $\mathbb{N}$. If we replace $\mathbb{C}$ by $\mathbb{Q}$, this result is no more true in general (see Example 3.5.23).

**Remark 3.3.25.** — Several authors proved "uniform" Strong Artin approximation results, i.e. they proved the existence of a function $\beta$ satisfying Theorem 3.3.14 for a family of $(f_\lambda(y, z))_{\lambda \in \Lambda}$ which satisfy tameness properties that we do not describe here. The main example is Theorem 3.3.12 that asserts that the Artin functions of polynomials in $n + m$ variables of degree less than $d$ are uniformly bounded. There are also two types of proof for these kind of "uniform" Strong Artin approximation results : the ones using ultraproducts (see Theorem 4.2 of [BDLvdD79] which is a generalization of Theorem 3.3.12 where the base field is not fixed, or Theorems 8.2 and 8.4 of [DeLi80] where uniform Strong Artin approximation results are proven for families of polynomials whose coefficients depend analytically on some parameters) and the ones using the scheme of proof due to Artin (see [ElTo96] where more or less the same results as those of [BDLvdD79] and [DeLi80] are proven).

3.3.4. Effectivity of the behaviour of Artin functions: some examples. — In general the proofs of Strong Artin Approximation results do not give much information about the Artin functions, since ultraproducts methods use a proof by contradiction (see also Remark 3.3.13). The problem of finding estimates of Artin functions was raised first in [Ar70] and very few general results are known (the only ones in the case of Greenberg Theorem are Theorems 3.3.8, 3.3.9 and Remark 3.3.10 and Remark 3.3.13 in the general case). We
give here a list of examples for which we can give non trivial effective behaviour about their Artin function.

3.3.4.1. Artin-Rees Lemma. — The following result has been known for long by the specialists and has been communicated to the author by M. Hickel:

**Theorem 3.3.26.** — Let \( f(y) \in A[y]^r \) be a vector of linear polynomials with coefficients in a Noetherian ring \( A \). Let \( I \) be an ideal of \( A \). Then there exists a constant \( c_0 \geq 0 \) such that:

\[
\forall c \in \forall \bar{y} \in A^m \text{ such that } f(\bar{y}) \in I^{c+c_0}
\]

\[
\exists \bar{y} \in A^m \text{ such that } f(\bar{y}) = 0 \text{ and } \bar{y} - \bar{y} \in I^c.
\]

This theorem asserts that the function \( \beta \) of Theorem 3.3.16 is bounded by the function \( c \mapsto c + c_0 \). Moreover let us remark that this theorem is valid for general Noetherian ring and general ideals \( I \) if \( A \). This can be compared with the fact that, for linear equations, Theorem 3.2.16 is true for any Noetherian ring \( A \) without Henselian condition (see Remark 3.2.20).

**Proof.** — For convenience, let us assume that there is only one linear polynomial:

\[
f(y) = a_1y_1 + \cdots + a_m y_m.
\]

Let us denote by \( I \) the ideal of \( A \) generated by \( a_1, \ldots, a_m \). Artin-Rees Lemma implies that there exists \( c_0 > 0 \) such that \( I \cap I^{c+c_0} \subset I.I^c \) for any \( c \geq 0 \). If \( \bar{y} \in A^m \) is such that \( f(\bar{y}) \in I^{c+c_0} \) then, since \( f(\bar{y}) \in I \), there exists \( \varepsilon \in (I^c)^m \) such that \( f(\bar{y}) = f(\varepsilon) \). If we define \( \bar{y}_i := \bar{y}_i - \varepsilon_i \), for \( 1 \leq i \leq m \), we have the result. \( \square \)

We have the following result whose proof is similar:

**Proposition 3.3.27.** — Let \((A, m_A)\) be a Henselian excellent local ring, \( I \) an ideal of \( A \) generated by \( a_1, \ldots, a_q \) and \( f(y) \in A[y]^r \). Set

\[
F_i(y, z) := f_i(y) + a_1 z_{i,1} + \cdots + a_q z_{i,q} \in A[y, z], \quad 1 \leq i \leq r
\]

where the \( z_{i,k} \)'s are new variables and let \( F(y, z) \) be the vector whose coordinates are the \( F_i(y, z) \)'s. Let us denote by \( \beta \) the Artin function of \( f(y) \) seen as a vector of polynomials of \( A[y] \) and \( \gamma \) the Artin function of \( F(y, z) \in A[y, z]^r \). Then there exists a constant \( c_0 \) such that:

\[
\beta(c) \leq \gamma(c) \leq \beta(c + c_0), \quad \forall c \in \mathbb{N}.
\]

**Proof.** — Let \( \bar{y} \in A^m \) satisfies \( f(\bar{y}) \in m_A^{\gamma(c)} A \). Then there exists \( \bar{z} \in A^{qr} \) such that \( F(\bar{y}, \bar{z}) \in m_A^{\gamma(c)} \) (we still denote by \( \bar{y} \) a lifting of \( \bar{y} \) in \( A^m \)). Thus there exists \( \tilde{y} \in A^m \) and \( \tilde{z} \in A^{qr} \) such that \( F(\tilde{y}, \tilde{z}) = 0 \) and \( \tilde{y} - \bar{y}, \tilde{z} - \bar{z} \in m_A^c \). Thus \( f(\tilde{y}) = 0 \) in \( A \).
Let \( c_0 \) be a constant such that \( I \cap m_A^{c+c_0} \subset I \cdot m_A^c \) for all \( c \in \mathbb{N} \) (such constant exists by the Artin-Rees Lemma). Let \( \overline{y} \in A^m, \overline{z} \in A^n \) satisfy \( F(\overline{y}, \overline{z}) \in m_A^{\beta(c+c_0)} \). Then \( f(\overline{y}) \in m_A^{\beta(c+c_0)} + I \). Thus there exists \( \tilde{y} \in A^m \) such that \( f(\tilde{y}) \in I \) and \( \tilde{y} - y \in m_A^{c+c_0} \). Thus \( F(\tilde{y}, \overline{z}) \in m_A^{c+c_0} \cap I \). Then we conclude by following the proof of Theorem 3.3.26.

**Remark 3.3.28.** — By Theorem [3.2.17] in order to study the behaviour of the Artin function of some ideal we may assume that \( A \) is a complete local ring. Let us assume that \( A \) is an equicharacteristic local ring. Then \( A \) is the quotient of a power series ring over a field by Cohen Structure Theorem [Mat80]. Thus Proposition 3.3.27 allows us to reduce the problem to the case \( A = \mathbb{k}[x_1, \ldots, x_n] \) where \( \mathbb{k} \) is a field.

### 3.3.4.2. Izumi’s Theorem and Diophantine Approximation

Let \((A, m_A)\) be a Noetherian local ring. We denote by \( \nu \) the \( m_A \)-adic order on \( A \), i.e.

\[
\nu(x) := \max\{n \in \mathbb{N} \mid x \in m_A^n\}
\]

for any \( x \neq 0 \).

We always have \( \nu(x) + \nu(y) \leq \nu(xy) \) for all \( x, y \in A \). But we do not have the equality in general. For instance, if \( A := \mathbb{C}[x,y]/(x^2 - y^3) \) then \( \nu(x) = \nu(y) = 1 \) but \( \nu(x^2) = \nu(y^3) = 3 \). Nevertheless we have the following theorem:

**Theorem 3.3.29 (Izumi’s Theorem).** — [Iz85][Rec89] Let \((A, m_A)\) be a local Noetherian ring. Let us assume that \( A \) is analytically irreducible, i.e. \( \hat{A} \) is irreducible. Then there exist \( b \geq 1 \), and \( d \geq 0 \) such that

\[
\forall x, y \in A, \quad \nu(xy) \leq b(\nu(x) + \nu(y)) + d.
\]

This result implies easily the following corollary using Corollary 3.3.27.

**Corollary 3.3.30.** — [Iz95][Ron06a] Let us consider the polynomial

\[
f(y) := y_1y_2 + a_3y_3 + \cdots + a_my_m,
\]

with \( a_3, \ldots, a_m \in A \) where \((A, m_A)\) is a Noetherian local ring such that \( (a_3, \ldots, a_m) \) is analytically irreducible. Then there exist \( b \geq 1 \) and \( d \geq 0 \) such that the Artin function \( \beta \) of Theorem 3.3.16 satisfies \( \beta(c) \leq bc + d \) for all \( c \in \mathbb{N} \).

**Proof.** — By Proposition 3.3.27 we have to prove that the Artin function \( \beta \) of \( y_1y_2 \in A[y] \) is bounded by an affine function if \( A \) is analytically irreducible. Thus let \( y_1, y_2 \in A \) satisfy \( y_1y_2 \in m_A^{2bc+d} \) where \( b \) and \( d \) satisfies Theorem 3.3.29. This means that

\[
2bc + d \leq \nu(y_1y_2) \leq b(\nu(y_1) + \nu(y_2)) + d.
\]
Thus $\nu(\tilde{y}_1) \geq c$ or $\nu(\tilde{y}_2) \geq c$. In the first case we denote $\tilde{y}_1 = 0$ and $\tilde{y}_2 = \tilde{y}_2$ and in the second case we denote $\tilde{y}_1 = \tilde{y}_1$ and $\tilde{y}_2 = 0$. Then $\tilde{y}_1 \tilde{y}_2 = 0$ and $\tilde{y}_1 - \tilde{y}_1, \tilde{y}_2 - \tilde{y}_2 \in m_A$. 

**Hints on the proof of Theorem 3.3.29 in the complex analytic case**

According to the theory of Rees valuations, there exists discrete valuations $\nu_1, \ldots, \nu_k$ such that $\nu(x) = \min\{\nu_1(x), \ldots, \nu_k(x)\}$ (they are called the Rees valuation of $\nu$). The valuation rings associated to $\nu_1, \ldots, \nu_k$ are the valuation rings associated to the irreducible components of the exceptional divisor of the normalized blowup of $m_A$.

Since $\nu_i(xy) = \nu_i(x) + \nu_i(y)$ for any $i$, in order to prove the theorem we have to see that there exists $a \geq 1$ such that $\nu_i(x) \leq ax_j(x)$ for any $x \in A$ and any $i$ and $j$. If $A$ is a complex analytic local ring, following S. Izumi’s proof, we may reduce the problem to the case $\dim(A) = 2$ by using a Bertini type theorem, and then assume that $A$ is normal by using an inequality on the reduced order proved by D. Rees. Then let us consider a resolution of singularities of Spec($A$) (denoted by $\pi$) that factors through the normalized blow-up of $m_A$. In this case, let us denote by $E_1, \ldots, E_s$ the irreducible components of the exceptional divisor of $\pi$. Let us denote $e_{i,j} := E_i.E_j$ for all $1 \leq i, j \leq s$. Let $x$ be an element of $A$. This element defines a germ of analytic hypersurface whose total transform $T_x$ may be written $T_x = S_x + \sum_{j=1}^s m_j E_j$ where $S_x$ is the strict transform of $\{x = 0\}$ and $m_i = \nu_i(x), 1 \leq i \leq s$. Then we have

$$0 = T_x.E_i = S_x.E_i + \sum_{j=1}^s m_j e_{i,j}.$$ 

Since $S_x.E_i \geq 0$ for any $i$, the vector $(m_1, \ldots, m_s)$ is contained in the convex cone $C$ defined by $x_i \geq 0, 1 \leq i \leq s$, and $\sum_{j=1}^s e_{i,j} x_j \leq 0, 1 \leq i \leq s$. To prove the theorem, it is enough to prove that $C$ is included in $x_i > 0, 1 \leq i \leq s$. Let assume that it is not the case. Then, after renumbering the $E_i$’s, we may assume that $(x_1, \ldots, x_1, 0, \ldots, 0) \in C$ where $x_1 > 0, 1 \leq i < l < s$. Since $e_{i,j} \geq 0$ for all $i \neq j$, $\sum_{j=1}^s e_{i,j} x_j = 0$ for $l < i \leq s$ implies that $e_{i,j} = 0$ for all $l < i \leq s$ and $1 \leq j \leq l$. This contradicts the fact that the exceptional divisor of $\pi$ is connected (since $A$ is an integral domain).

Let us mention that Izumi’s Theorem is the key ingredient to prove the following result:

**Corollary 3.3.31.** — Let $(A, m_A)$ be a regular excellent Henselian domain. Let us denote by $\mathbb{K}$ and $\hat{\mathbb{K}}$ the fraction fields of $A$
and \( \hat{A} \) respectively. Let \( z \in \hat{K} \setminus K \) be algebraic over \( K \). Then
\[
\exists a \geq 1, C \geq 0, \forall x \in A \forall y \in A^* \left| z - \frac{x}{y} \right| \geq C|y|^a
\]
where \(|u| := e^{-\nu(a)}\) and \( \nu \) is the usual \( m_A \)-adic valuation.

This result is equivalent to the following:

**Corollary 3.3.32.** — [Ron06b][Hic08][ItIz08] Let \((A, m_A)\) be an excellent Henselian local ring and let \( f_1(y_1, y_2), \ldots, f_r(y_1, y_2) \in A[y_1, y_2] \) be homogeneous polynomials. Then the Artin function of \( f_1, \ldots, f_r \) is bounded by an affine function.

### 3.3.4.3. Reduction to one quadratic equation and examples

In general Artin functions are not bounded by affine functions as in Theorem [3.3.1]. Here is such an example:

**Example 3.3.33.** — [Ron05b] Set \( f(y_1, y_2, y_3) := y_1^2 - y_2^2 y_3 \in k[[x_1, x_2]][y_1, y_2, y_3] \) where \( k \) is a field of characteristic zero. Let us denote by \( h(T) := \sum_{i=1}^{\infty} a_i T^i \in \mathbb{Q}[T] \) the power series such that \((1 + h(T))^2 = 1 + T\). Let us denote
\[
y_1^{(c)} := x_1^{2c+2} \left( 1 + \sum_{i=1}^{c+1} a_i \frac{x_1^{2(i-1)}}{x_1^i} \right) = x_1^{2c+2} + \sum_{i=1}^{c+1} a_i x_1^{2(c-i+1)} x_2^i,
y_2^{(c)} := x_1^{2c+1},
y_3^{(c)} := x_1^2 + x_2.
\]

Then in the ring \( k[\{x_i\}][x_1] \) we have
\[
f(y_1^{(c)}, y_2^{(c)}, y_3^{(c)}) = \left( \frac{y_1^{(c)}}{y_2^{(c)}} - y_3^{(c)} \right) y_2^{(c)} = \left( \frac{y_1^{(c)}}{y_2^{(c)}} - x_1 \left( 1 + \frac{x_2}{x_1} \right) \right) y_2^{(c)}
\]
\[
= \frac{y_1^{(c)}}{y_2^{(c)}} - x_1 \left( 1 + h \left( \frac{x_2}{x_1} \right) \right) \left( \frac{y_1^{(c)}}{y_2^{(c)}} + x_1 \left( 1 + h \left( \frac{x_2}{x_1} \right) \right) \right) y_2^{(c)}.
\]

Thus we see that \( f(y_1^{(c)}, y_2^{(c)}, y_3^{(c)}) \in (x)^{2c+4c} \) for all \( c \geq 2 \). But if \( (y_1^{(c)}, y_2^{(c)}, y_3^{(c)}) \in k[[x_1, x_2]]^3 \) is a solution of \( f = 0 \) then

1) Either \( y_3 \) is a square in \( k[[x_1, x_2]] \). But \( \sup_{z \in k[x]} \ord(y_3^{(c)} - z^2) = c \).

2) Either \( y_3 \) is not a square, hence \( \tilde{y}_1 = \tilde{y}_2 = 0 \). But \( \ord(y_1^{(c)}) - 1 = \ord(y_2^{(c)}) = 2c + 1 \).
In any case we have
\[
\sup(\min\{\ord(y_1^{(c)} - \tilde{y}_1), \ord(y_2^{(c)} - \tilde{y}_2), \ord(y_3^{(c)} - \tilde{y}_3)\}) \leq 2c + 1.
\]
This proves that the Artin function \( f \) is bounded from below by a polynomial function of degree 2. Thus Theorem 3.3.1 does not extend to \( k[x_1, ..., x_n] \) if \( n \geq 2 \).

In [Ron06a] another example is given: the Artin function of the polynomial \( y_1 y_2 - y_3 y_4 \in k[x_1, x_2, x_3][y_1, y_2, y_3, y_4] \) is bounded from below by a polynomial function of degree 2. Both examples are the only known examples of Artin functions which are not bounded by an affine function.

In general, in order to investigate bounds on the growth of Artin functions in general, we can reduce the problem as follows, using a trick of [Ron10b].

**Lemma 3.3.34. — [Be77b]** For any \( f(y) \in A[y]^r \) or \( A[y]^r \) the Artin function of \( f \) is bounded by the Artin function of
\[
g(y) := f_1(y)^2 + y_1(f_2(y)^2 + y_1(f_3(y)^2 + \cdots)^2).
\]

**Proof.** — Indeed, if \( \beta \) is the Artin function of \( g \) and if \( f(\tilde{y}) \in m_A^{\beta(c)} \) then \( g(\tilde{y}) \in m_A^{\beta(c)} \). Thus there exists \( \tilde{y} \in A^n \) such that \( g(\tilde{y}) = 0 \) and \( \tilde{y}_i - y_i \in m_A^c \). But clearly \( g(\tilde{y}) = 0 \) if and only if \( f(\tilde{y}) = 0 \). This proves the lemma.

This allows us to assume that \( r = 1 \) and we define \( f(y) := f_1(y) \). If \( f(y) \) is not irreducible, then we may write \( f = h_1 \cdots h_s \), where \( h_i \in A[y] \) is irreducible for \( 1 \leq i \leq s \), and the Artin function of \( f \) is bounded by the sum of the Artin functions of the \( h_i \)’s. Hence we may assume that \( f(y) \) is irreducible.

We have the following lemma:

**Lemma 3.3.35. — [Ron09b] [Ron10b]** For any \( f(y) \in A[y] \), where \( A \) is a complete local ring, the Artin function of \( f(y) \) is bounded by the Artin function of the polynomial
\[
P(u, x, z) := f(y)u + x_1 z_1 + \cdots + x_m z_m \in B[x, z, u]
\]
where \( B := A[y] \).

**Proof.** — Let us assume that \( f(\tilde{y}) \in m_A^{\beta(c)} \) where \( \beta \) is the Artin function of \( P \). By replacing \( f(y) \) by \( f(y^0 + y) \), where \( y^0 \in A \) is such that \( y^0_i - \tilde{y}_i \in m_A^c \).
$1 \leq i \leq m$, we may assume that $\overline{y}_i \in m_A$ for $1 \leq i \leq m$. Then there exists $z_i(y) \in A[y]$, $1 \leq i \leq m$, such that

$$f(y) + \sum_{i=1}^m (y_i - \overline{y}_i)z_i(y) \in (m_A + (y))^\beta(c).$$

Thus there exists $u(y), f_i(y), z_i(y) \in A[y]$, $1 \leq i \leq m$, such that

$$u(y) - 1, z_i(y) - z_i(y), x_i(y) - (y_i - \overline{y}_i) \in (m_A + (y))\gamma, 1 \leq i \leq n$$

and $f(y)u(y) + \sum_{i=1}^m x_i(y)z_i(y) = 0$.

In particular $u(y)$ is invertible in $A[y]$ if $c > 0$. Let us assume that $c \geq 2$. In this case the matrix of the partial derivatives of $(x_i(y), 1 \leq i \leq m)$ with respect to $y_1, \ldots, y_m$ has determinant equal to 1 modulo $m_A + (y)$. By the Henselian property there exist $y_{i,c} \in m_A$ such that $x_i(y_{1,c}, \ldots, y_{m,c}) = 0$ for $1 \leq i \leq m$. Hence, since $u(y_{i,c})$ is invertible, $f(y_{i,c}, \ldots, y_{m,c}) = 0$ and $y_{i,c} - \overline{y}_i \in m_A^\epsilon$, $1 \leq i \leq m$.

Thus, by Corollary 3.3.27 in order to study the general growth of Artin functions, it is enough to study the Artin function of the polynomial $y_1y_2 + y_3y_4 + \cdots + y_{2m+1}y_{2m} \in A[y]$ where $A$ is a complete local ring.

### 3.4. Examples of Applications

In this part we give some examples of applications of Theorem 3.2.17 and Corollary 3.3.16.

**Proposition 3.4.1.** — Let $A$ be an excellent Henselian local ring. Then $A$ is reduced (resp. is an integral domain, resp. an integrally closed domain) if and only if $\hat{A}$ is reduced (resp. is an integral domain, resp. an integrally closed domain).

**Proof.** — If $\hat{A}$ is not reduced, then there exists $\hat{y} \in \hat{A}$, $\hat{y} \neq 0$, such that $\hat{y}^k = 0$ for some positive integer $k$. Thus we apply Theorem 3.2.17 to the polynomial $y^k$ with $c \geq \text{ord}(\hat{y}) + 1$ in order to find $\hat{y} \in A$ such that $\hat{y}^k = 0$ and $\hat{y} \neq 0$.

In order to prove that $\hat{A}$ is an integral domain if $A$ is an integral domain, we apply the same procedure to the polynomial $y_1y_2$.

If $A$ is an integrally closed domain, then $\hat{A}$ is an integral domain. Let $P(z) := z^d + \hat{a}_1z^{d-1} + \cdots + \hat{a}_d \in \hat{A}[z]$, $\hat{f}, \hat{g} \in \hat{A}$, $\hat{g} \neq 0$, satisfy $P\left(\frac{\hat{f}}{\hat{g}}\right) = 0$, i.e.

$$\hat{f}^d + \hat{a}_1\hat{f}^{d-1}\hat{g} + \cdots + \hat{a}_d\hat{g}^d = 0.$$ 

By Theorem 3.2.17 for any $c \in \mathbb{N}$, there...
exist $\tilde{a}_{i,c}, \tilde{f}_c, \tilde{g}_c \in A$ such that $\tilde{f}_c^d + \tilde{a}_{1,c, \tilde{f}_c} + \tilde{\cdots} + \tilde{a}_{d,c, \tilde{f}_c} = 0$ and $\tilde{f}_c - \tilde{f}, \tilde{g}_c - \tilde{g} \in \mathfrak{m}_A^c$. Then for $c > c_0$, for some integer $c_0$, $\tilde{g}_c \neq 0$. Since $A$ is an integrally closed domain, then $\tilde{f}_c \in (\tilde{g}_c)$ for $c > c_0$. Thus $\tilde{f} \in (\tilde{g}) + \mathfrak{m}_A^c$ for $c$ large enough. By Nakayama’s Lemma this implies that $\tilde{f} \in (\tilde{g})$ and $\hat{A}$ is integrally closed.

**Proposition 3.4.2.** — Let $A$ be an excellent Henselian local ring. Let $Q$ be a primary ideal of $A$. Then $QA$ is a primary ideal of $\hat{A}$.

**Proof.** — Let $\hat{f} \in \hat{A}$ and $\hat{g} \in \hat{A} \setminus \sqrt{Q \hat{A}}$ satisfy $\hat{f}\hat{g} \in Q \hat{A}$. By Theorem 3.2.17 for any $c \in \mathbb{N}$, there exist $\tilde{f}_c, \tilde{g}_c \in A$ such that $f_c g_c \in Q$ and $\tilde{f}_c - \hat{f}, \tilde{g}_c - \hat{g} \in \mathfrak{m}_A^c$. For some $c$ large enough, $\tilde{g}_c \notin \sqrt{Q}$. Since $A$ is a primary ideal, this proves that $\tilde{f}_c \in Q$ for $c$ large enough, hence $\hat{f} \in Q \hat{A}$.

**Corollary 3.4.3.** — Let $A$ be an excellent Henselian local ring. Let $I$ be an ideal of $A$ and let $I = Q_1 \cap \cdots \cap Q_s$ be a primary decomposition of $I$ in $A$. Then $Q_1 \hat{A} \cap \cdots \cap Q_s \hat{A}$ is a primary decomposition of $I\hat{A}$.

**Proof.** — Since $I = \bigcap_{i=1}^s Q_i$, then $I \hat{A} = \bigcap_{i=1}^s (Q_i \hat{A})$ by faithfull flatness (or by Theorem 3.2.17 for linear equations). We conclude with the help of Proposition 3.4.2.

**Corollary 3.4.4.** — [Iz92] Let $A$ be an excellent Henselian local integrally closed domain. If $\hat{f} \in \hat{A}$ and if there exists $\hat{g} \in \hat{A}$ such that $\hat{f}\hat{g} \in A \setminus \{0\}$, then there exists a unit $\hat{u} \in \hat{A}$ such that $\hat{u}\hat{f} \in A$.

**Proof.** — Let $(\hat{f}\hat{g})A = Q_1 \cap \cdots \cap Q_s$ be a primary decomposition of the principal ideal of $A$ generated by $\hat{f}\hat{g}$. Since $A$ is an integrally closed domain, it is a Krull ring and $Q_i = p_i^{(n_i)}$ for some prime ideal $p_i$, $1 \leq i \leq s$, where $p_i^{(n)}$ denote the $n$-th symbolic power of $p$ (see [Mat80] p.88). In fact $n_i := \nu_{p_i}(\hat{f}\hat{g})$ where $\nu_{p_i}$ is the $p_i$-adic valuation of the valuation ring $A_{p_i}$. By Corollary 3.4.3 $p_1^{(n_1)} \hat{A} \cap \cdots \cap p_s^{(n_s)} \hat{A}$ is a primary decomposition of $(\hat{f}\hat{g})\hat{A}$. Since $\nu_{p_i}$ are valuations, then

$$\hat{f}\hat{A} = p_1^{(k_1)} \hat{A} \cap \cdots \cap p_s^{(k_s)} \hat{A} = \left( p_1^{(k_1)} \cap \cdots \cap p_s^{(k_s)} \right) \hat{A}$$

for some non negative integers $k_1, \ldots, k_s$. Let $h_1, \ldots, h_r \in A$ be generators of $p_1^{(k_1)} \cap \cdots \cap p_s^{(k_s)}$. Then $\hat{f} = \sum_{i=1}^r \hat{a}_i h_i$ and $h_i = \hat{b}_i \hat{f}$ for some $\hat{a}_i, \hat{b}_i \in A$, which completes the proof.
1 \leq i \leq r. Thus \( \sum_{i=1}^{r} \widehat{a}_i \widehat{b}_i = 1 \), since \( \widehat{A} \) is an integral domain. Thus one of the \( \widehat{b}_i \)'s is invertible and we choose \( \widehat{u} \) to be this invertible \( \widehat{b}_i \).

**Corollary 3.4.5.** — [K072] Let \( A \) be an excellent Henselian local ring. For \( f(y) \in A[y]^r \) let \( I \) be the ideal of \( A[y] \) generated by \( f_1(y), \ldots, f_r(y) \). Let us assume that \( \text{ht}(I) = m \). Let \( \widehat{y} \in \widehat{A}^m \) satisfy \( f(\widehat{y}) = 0 \). Then \( \widehat{y} \in \widehat{A}^m \).

**Proof.** — Set \( \mathfrak{p} := (y_1 - \widehat{y}_1, \ldots, y_m - \widehat{y}_m) \). It is a prime ideal of \( \widehat{A} \) and \( \text{ht}(\mathfrak{p}) = m \). Of course \( I \widehat{A} \subset \mathfrak{p} \) and \( \text{ht}(I \widehat{A}) = m \) by Corollary 3.4.3. Thus \( \mathfrak{p} \) is of the form \( \mathfrak{p}' \widehat{A} \) where \( \mathfrak{p}' \) is minimal prime of \( I \). Then \( \widehat{y} \in \widehat{A}^m \) is the only common zero of all the elements of \( \mathfrak{p}' \). By Theorem 3.2.17, \( \widehat{y} \) can be approximated by a common zero of all the elements of \( \mathfrak{p}' \) which is in \( \widehat{A}^m \). Thus \( \widehat{y} \in \widehat{A}^m \).

**Proposition 3.4.6.** — [KPPRM78, Po86] Let \( A \) be an excellent Henselian local ring. Then \( A \) is a unique factorization domain if and only if \( \widehat{A} \) is a unique factorization domain.

**Proof.** — If \( \widehat{A} \) is a unique factorization domain, then any irreducible element of \( \widehat{A} \) is prime. Thus any irreducible element of \( A \) is prime. Since \( A \) is a Noetherian integral domain, it is a unique factorization domain.

Let us assume that \( \widehat{A} \) is a Noetherian integral domain but not a unique factorization domain. Thus there exists an irreducible element \( \widehat{x}_1 \in \widehat{A} \) that is not prime. This equivalent to

\[
\exists \widehat{x}_2, \widehat{x}_3, \widehat{x}_4 \in \widehat{A} \text{ such that } \widehat{x}_1 \widehat{x}_2 - \widehat{x}_3 \widehat{x}_4 = 0
\]

\[
\beta \widehat{x}_1, \widehat{x}_2 \in \widehat{A} \text{ such that } \widehat{x}_1 \widehat{x}_1 - \widehat{x}_3 = 0 \text{ and } \widehat{x}_2 \widehat{x}_2 - \widehat{x}_4 = 0
\]

and \( \beta \widehat{y}_1, \widehat{y}_2 \in m_A \widehat{A} \) such that \( \widehat{y}_1 \widehat{y}_2 - \widehat{x}_1 = 0 \).

Let us denote by \( \beta \) the Artin function of

\[
f(y, z) := (\widehat{x}_1 y_1 - \widehat{x}_3, \widehat{x}_2 z_2 - \widehat{x}_4, y_1 y_2 - \widehat{x}_1) \in \widehat{A}[y][z].
\]

Since \( f(y, z) \) has no solution in \( (m_A \widehat{A})^2 \times \widehat{A}^2 \), by Remark 3.3.18 \( \beta \) is a constant, and \( f(y, z) \) has no solution in \((m_A \widehat{A})^2 \times \widehat{A}^2\) modulo \( m_A^3 \).

On the other hand by Theorem 3.2.17 applied to \( x_1 x_2 - x_3 x_4 \), there exists \( x_i \in A \), \( 1 \leq i \leq 4 \), such that \( \widehat{x}_1 \widehat{x}_2 - \widehat{x}_3 \widehat{x}_4 = 0 \) and \( \widehat{x}_i - \widehat{x}_i \in m_A^{3+1}, 1 \leq i \leq 4 \).

Hence

\[
g(y, z) := (\widehat{x}_1 y_1 - \widehat{x}_3, \widehat{x}_2 z_2 - \widehat{x}_4, y_1 y_2 - \widehat{x}_1) \in \widehat{A}[y][z]
\]

has no solution in \((m_A \widehat{A})^2 \times \widehat{A}^2\) modulo \( m_A^3 \), hence has no solution in \((m_A A)^2 \times A^2\). This means that \( \widehat{x}_1 \) is an irreducible element of \( A \) but it is not prime. Hence \( A \) is not a unique factorization domain. \( \square \)
3.5. Approximation with constraints

We will now discuss the problem of the Artin Approximation with constraints that is the following:

Problem 1 (Artin Approximation with constraints):
Let \( A \) be an excellent Henselian local subring of \( k[[x_1, \ldots, x_n]] \) and \( f(y) \in A[y]^r \).
Let us assume that we have a formal solution \( \hat{y} \in \hat{A}^m \) of \( f = 0 \) and assume moreover that
\[
\hat{y}_i(x) \in \hat{A} \bigcap k[[x_j, j \in J_i]]
\]
for some subset \( J_i \subset \{1, \ldots, n\} \), \( 1 \leq i \leq m \).
Is it possible to approximate \( \hat{y}(x) \) by a solution \( \tilde{y}(x) \in A^m \) of \( f = 0 \) such that
\[
\tilde{y}_i(x) \in k[[x_j, j \in J_i]] \text{, } 1 \leq i \leq m?
\]

Another problem is the following:

Problem 2 (Strong Artin Approximation with constraints):
Let us consider \( f(y) \in k[[x]]^r \) and \( J_i \subset \{1, \ldots, n\} \), \( 1 \leq i \leq m \). Does there exist a function \( \beta : \mathbb{N} \rightarrow \mathbb{N} \) such that:
for all \( c \in \mathbb{N} \) and all \( \overline{y}_i(x) \in k[[x_j, j \in J_i]] \), \( 1 \leq i \leq m \), such that
\[
f(\overline{y}(x)) \in (x)^{\beta(c)},
\]
there exist \( \tilde{y}_i(x) \in k[[x_j, j \in J_i]] \) such that \( f(\overline{y}(x)) = 0 \) and \( \tilde{y}_i(x) - \overline{y}_i(x) \in (x)^c \), \( 1 \leq i \leq m \)?

If such function \( \beta \) exists, the smallest function satisfying this property is called the Artin function of the system \( f = 0 \).

Let us remark that we have already given a positive answer to a similar weaker problem (see Proposition 3.3.24). The answer will be no in general for both problems and yes for some particular cases. We present here the positive and negative results concerning these problems. We will see that some systems yield a positive answer to Problem 2 but a negative answer to Problem 1.

3.5.1. Examples. — First of all we give here a list of examples that show that there is no hope, in general, to have a positive answer to Problem 1 without any more specific hypothesis, even if \( A \) is the ring of algebraic or convergent power series. These examples are constructed by looking at the Artin Approximation Problem for equations involving differentials (Examples 3.5.3 and 3.5.6) and operators on germs of functions (Examples 3.5.4 and 3.5.5). To construct these examples, the following lemma will be used repeatedly:
Lemma 3.5.1. — Let \((A, m_A)\) be a Noetherian local ring and let \(B\) be a Noetherian local subring of \(\hat{A}[y]\) such that \(\hat{B} = \hat{A}[y]\). For any \(P(y) \in B\) and \(\hat{y} \in (m_A.A)^m\), \(P(\hat{y}) = 0\) if and only if there exists \(\hat{h}(y) \in B^m\) such that

\[
P(y) + \sum_{i=1}^{m} (y_i - \hat{y}_i) \hat{h}(y) = 0.
\]

Proof of Lemma 3.5.1. — If \(P(\hat{y}) = 0\) then, by Taylor expansion, we have:

\[
P(y) - P(\hat{y}) = \sum_{\alpha \in \mathbb{N}^m \setminus \{0\}} \frac{1}{\alpha_1! \cdots \alpha_m!} (y_1 - \hat{y}_1)^{\alpha_1} \cdots (y_m - \hat{y}_m)^{\alpha_m} \frac{\partial^\alpha P(\hat{y})}{\partial y^\alpha}.
\]

Thus there exists \(\hat{h}(y) \in A[y]^m\) such that

\[
P(y) + \sum_{i=1}^{m} (y_i - \hat{y}_i) \hat{h}(y) = 0.
\]

Since \(B \longrightarrow \hat{B} = \hat{A}[x]\) is faithfully flat and we may assume that \(\hat{h}(y) \in B\) (See Example 3.1.4).

On the other hand if \(P(y) + \sum_{i=1}^{m} (y_i - \hat{y}_i) \hat{h}(y) = 0\), by substitution of \(y_i\) by \(\hat{y}_i\), we get \(P(\hat{y}) = 0\).

Example 3.5.2. — Let us consider \(P(x, y, z) \in k[[x, y, z]]\) where \(x, y\) and \(z\) are single variables and \(\hat{y} \in (x).k[[x]]\). Then \(P(x, \hat{y}, \frac{\partial \hat{y}}{\partial x}) = 0\) if and only if \(P(x, \hat{y}, \hat{z}) = 0\) and \(\hat{z} - \frac{\partial \hat{y}}{\partial x} = 0\).

Moreover \(\hat{z} - \frac{\partial \hat{y}}{\partial x} = 0\) if and only if \(\hat{z} - \left(\frac{\hat{y}(x+t) - \hat{y}(x)}{t}\right) \in (t)k[[x, t]]\). By Lemma 3.5.1 this is equivalent to: there exist \(\hat{h}(x, t, u), \hat{k}(x, t, u) \in k[[x, t, u]]\) such that

\[
t\hat{z}(x) - \hat{y}(u) - \hat{y}(x) + t^2 \hat{h}(x, t) + (u - x - t) \hat{k}(x, t, u) = 0.
\]

Finally we see that

\[
P(x, \hat{y}(x), \frac{\partial \hat{y}}{\partial x}(x)) = 0 \iff \\
\exists \hat{z}(x) \in k[[x]], \hat{h}(x, t, u), \hat{k}(x, t, u), \hat{l}(x, t, u) \in k[[x, t, u]], \hat{y}(u) \in k[[u]]\ s.t.
\]

\[
\begin{cases}
P(x, \hat{y}(x), \hat{z}(x)) = 0 \\
t\hat{z}(x) - \hat{y}(u) - \hat{y}(x) + t^2 \hat{h}(x, t, u) + (u - x - t) \hat{k}(x, t, u) = 0 \\
\hat{y}(u) - \hat{y}(x) + (u - y) \hat{l}(x, t, u) = 0
\end{cases}
\]

Lemma 3.5.1 and Example 3.5.2 allow us to transform any system of equations involving partial differentials and compositions of power series into a system of algebraic equations whose solutions depend only on some of the \(x_i\)’s.
3.5. APPROXIMATION WITH CONSTRAINTS

Of course there exists plenty of examples of such systems of equations with algebraic or analytic coefficients that do not have algebraic or analytic solutions. These kinds of examples will provide counterexamples to Problem 1 as follows:

**Example 3.5.3.** — Let us consider the following differential equation: $y' = y$. The solutions of this equation are the convergent power series $ce^x \in \mathbb{C}\{x\}$ where $c$ is a complex number.

On the other hand, by Example 3.5.2, $\hat{y}(x)$ is convergent power series solution of this equation if and only if there exists $\hat{y}_1(x_1) \in \mathbb{C}\{x_1\}$, $\hat{y}_2(x_2) \in \mathbb{C}\{x_2\}$ and $\hat{h}(x_1, x_2, x_3)$, $\hat{k}(x_1, x_2, x_3)$, $\hat{l}(x_1, x_2, x_3) \in \mathbb{C}\{x_1, x_2, x_3\}$ such that (with $\hat{y}_1 := \hat{y}$):

$$\begin{cases}
\hat{y}_2(x_2) - \hat{y}_1(x_1) = x_3\hat{y}_1(x_1) + x_2^2\hat{h}(x_1, x_2, x_3) + (x_2 - x_1 - x_3)\hat{k}(x_1, x_2, x_3) \\
\hat{y}_2(x_2) - \hat{y}_1(x_1) = (x_1 - x_2)\hat{l}(x_1, x_2, x_3)
\end{cases}$$

Thus the former system of equations has a convergent solution

$$(\hat{y}_1, \hat{y}_2, \hat{h}, \hat{k}, \hat{l}) \in \mathbb{C}\{x_1\} \times \mathbb{C}\{x_2\} \times \mathbb{C}\{x_1, x_2, x_3\},$$

but no algebraic solution in $\mathbb{C}\{x_1\} \times \mathbb{C}\{x_2\} \times \mathbb{C}\{x_1, x_2, x_3\}$.

**Example 3.5.4** (Kashiwara-Gabber’s Example). — ([Hir77], p. 75) Let us perform the division of $xy$ by

$$g := (x - y^2)(y - x^2) = xy - x^3 - y^3 + x^2y^2$$

as formal power series in $\mathbb{C}\{x, y\}$ with respect to the monomial $xy$ (see Example 3.1.10 in the introduction). The remainder of this division can be written $r(x) + s(y)$ where $r(x) \in (x)\mathbb{C}\{x\}$ and $s(y) \in (y)\mathbb{C}\{y\}$ since this remainder has no monomial divisible by $xy$. By symmetry, we get $r(x) = s(x)$, and by substituting $y$ by $x^2$ we get the following equation:

$$r(x^2) + r(x) - x^3 = 0.$$  

This relation yields the expansion

$$r(x) = \sum_{i=0}^{\infty} (-1)^i x^{3.2^i}$$

and shows that the remainder of the division is not algebraic. This proves that the equation

$$xy - gQ(x, y) - R(x) - S(y) = 0$$

has a convergent solution $(\hat{q}(x, y), \hat{r}(x), \hat{s}(y)) \in \mathbb{C}\{x, y\} \times \mathbb{C}\{x\} \times \mathbb{C}\{y\}$ but has no algebraic solution $(q(x, y), r(x), s(y)) \in \mathbb{C}\{x, y\} \times \mathbb{C}\{x\} \times \mathbb{C}\{y\}$.
Example 3.5.5 (Becker’s Example). — ([Be77b]) By direct computation we show that there exists a unique power series \( f(x) \in \mathbb{C}[x] \) such that \( f(x + x^2) = 2f(x) - x \) and that this power series is not convergent. But, by Lemma 3.5.1, we have:

\[
f(x + x^2) - 2f(x) + x = 0
\]

\[\iff \exists g(y) \in \mathbb{C}[y], h(x, y), k(x, y) \in \mathbb{C}[x, y] \text{ s.t.}
\]

\[
\begin{cases}
F_1 := g(y) - 2f(x) + x + (y - x - x^2)h(x, y) = 0 \\
F_2 := g(y) - f(x) + (x - y)k(x, y) = 0
\end{cases}
\]

Then this system of equations has solutions in \( \mathbb{C}[x] \times \mathbb{C}[y] \times \mathbb{C}[x, y]^2 \) but no solution in \( \mathbb{C}(x) \times \mathbb{C}(y) \times \mathbb{C}(x, y)^2 \), even no solution in \( \mathbb{C} \times \mathbb{C} \times \mathbb{C} \). All the examples of Section 3.5.1 involve components that depend on separate variables. Indeed, Example 3.5.2 shows that equations involving partial derivatives yield algebraic equations whose solutions have components with separate variables. In the case the variables are nested (i.e. \( y_i = y_i(x_1, \ldots, x_{s(i)}) \) for some integer \( i \), which is equivalent to say that \( J_i \) contains or is contained in \( J_j \) for any \( i \) and \( j \) with notations of Problems 1 and 2) it is not possible to construct a counterexample as we did in Section 3.5.1 from differential equations or equations as in Example 3.5.5. We will see, in the nested case, that the algebraic case is completely different from the analytic case. First of all in the algebraic case, we have a nested Artin approximation result as follows:

**Example 3.5.6.** — Set \( \hat{y}(x) := \sum_{i \geq 0} \frac{i!}{x^i+1} \in \mathbb{C}[x] \). This power series is divergent and we have shown in Example 3.1.15 that it is the only solution of the equation

\[x^2y' - y + x = 0 \text{ (Euler Equation).}\]

By Example 3.5.2, \( \hat{y}(x) \) is a solution of this differential equation if and only if there exist \( \hat{y}_2(x_2) \in \mathbb{C}[x_2] \) and \( \hat{k}(x_1, x_2, x_3), \hat{h}(x_1, x_2, x_3), \hat{l}(x_1, x_2, x_3) \in \mathbb{C}[x_1, x_2, x_3] \) such that \( (x := (x_1, x_2, x_3)):
\]

\[
\begin{cases}
x_1^2(\hat{y}_2(x_2) - \hat{y}_1(x_1)) - x_3\hat{y}_1(x_1) + x_3x_1 + x_3\hat{k}(x) + (x_1 + x_3 - x_2)\hat{h}(x) = 0 \\
\hat{y}_2(x_2) - \hat{y}_1(x_1) - (x_2 - x_1)\hat{l}(x) = 0
\end{cases}
\]

with \( \hat{y}_1(x_1) := \hat{y}(x_1) \). Thus this system has no solution in \( \mathbb{C}\{x_1\} \times \mathbb{C}\{x_2\} \times \mathbb{C}\{x_1, x_2, x_3\}^3 \) but it has solutions in \( \mathbb{C}[x_1] \times \mathbb{C}[x_2] \times \mathbb{C}[x_1, x_2, x_3]^3 \).

**Remark 3.5.7.** — By replacing \( f_1(y), \ldots, f_r(y) \) by \( g(y) := f_1(y)^2 + y_1(f_2(y)^2 + y_1(f_3(y)^2 + \cdots)^2)^2 \) in these examples as in the proof of Lemma 3.3.4 we can construct the same kind of examples involving only one equation. Indeed \( f_1 = f_2 = \cdots = f_r = 0 \) if and only if \( g = 0 \).

3.5.2. Nested Approximation in the algebraic case. — All the examples of Section 3.5.1 involve components that depends on separate variables. Indeed, Example 3.5.2 shows that equations involving partial derivatives yield algebraic equations whose solutions have components with separate variables. In the case the variables are nested (i.e. \( y_i = y_i(x_1, \ldots, x_{s(i)}) \) for some integer \( i \), which is equivalent to say that \( J_i \) contains or is contained in \( J_j \) for any \( i \) and \( j \) with notations of Problems 1 and 2) it is not possible to construct a counterexample as we did in Section 3.5.1 from differential equations or equations as in Example 3.5.5. We will see, in the nested case, that the algebraic case is completely different from the analytic case. First of all in the algebraic case, we have a nested Artin approximation result as follows:
Theorem 3.5.8 — [KPRM78] Let \((A, m_A)\) be an excellent Henselian local ring and \(f(x, y) \in A(x, y)^r\). Let \(\tilde{y}(x)\) be a solution of \(f = 0\) in \((m_A + (x))\hat{A}[x]_m\). Let us assume that \(\tilde{y}_i \in \hat{A}[x_1, \ldots, x_{s_i}]\), \(1 \leq i \leq m\), for integers \(s_i\), \(1 \leq s_i \leq n\). Then for any \(c \in \mathbb{N}\) there exists a solution \(\tilde{y}(x) \in A(x)^m\) such that \(\tilde{y}_i(x) \in A\langle x_1, \ldots, x_{s_i}\rangle\) and \(\tilde{y}(x) - \tilde{y}(x) \in (m_A + (x))^c\).

This result has a lot of applications and his one of the most important about Artin Approximation. The proof we present here uses the formalism of codes for algebraic power series and is a bit different from the classical one. The key point is the fact that \(A(x)\) satisfies Theorem 3.2.17 for any excellent Henselian local ring \(A\) (see Remark 3.2.22).

Proof of Theorem 3.5.8 — Lemma 3.5.9. — Let \(A\) be a complete normal local domain, \(u := (u_1, \ldots, u_n), v := (v_1, \ldots, v_m)\). Then
\[
A[u, v] = \{f \in A[u, v] / \exists s \in \mathbb{N}, g \in A(v, z_1, \ldots, z_s), \hat{z}_i \in (m_A + (u))A[u, v], f = g(v, \hat{z}_1, \ldots, \hat{z}_s)\}.
\]

Proof of Lemma 3.5.9. — Let us denote
\[
B := \{f \in A[u, v] / \exists s \in \mathbb{N}, g \in A(v, z_1, \ldots, z_s), \hat{z}_i \in (m_A + (u))A[u, v], f = g(v, \hat{z}_1, \ldots, \hat{z}_s)\}.
\]

Clearly \(B\) is a subring of \(A[u, v]\).
If \(f \in A[u, v]\) we can write \(f = f_0 + f_1\) where \(f_0 \in A\) and \(f_1 \in (m_A + (v))A(v)\). There exist \(F_1, \ldots, F_r \in A[u, v][X_1, \ldots, X_r]\) such that \(\frac{\partial(F_1, \ldots, F_r)}{\partial(X_1, \ldots, X_r)}\) is non-zero modulo \(m_A + (u, v, X)\) and such that the unique \((f_1, \ldots, f_r) \in (m_A + (u, v))A[u, v]^r\) with \(F(f_1, \ldots, f_r) = 0\) (by the Implicit Function Theorem) is such that \(f_1 = f_1\) (cf. Proposition 3.5.10). Let us write
\[
F_i := \sum_{\alpha, \beta} F_{i, \alpha, \beta} v^\alpha X^\beta, \quad 1 \leq i \leq r
\]
with \(F_{i, \alpha, \beta} \in A[u]\) for all \(i, \alpha, \beta\). We can write \(F_{i, \alpha, \beta} = F_{i, \alpha, \beta}^0 + \hat{z}_{i, \alpha, \beta}\) where \(F_{i, \alpha, \beta}^0 \in A\) and \(\hat{z}_{i, \alpha, \beta} \in (m_A + (u))A[u]\). Let us denote
\[
G_i := \sum_{\alpha, \beta} (F_{i, \alpha, \beta}^0 + \hat{z}_{i, \alpha, \beta}) v^\alpha X^\beta, \quad 1 \leq i \leq r
\]
where \(\hat{z}_{i, \alpha, \beta}\) are new variables. Let us denote by \(z\) the vector whose coordinates are the variables \(\hat{z}_{i, \alpha, \beta}\). Then \(\frac{\partial(G_1, \ldots, G_r)}{\partial(X_1, \ldots, X_r)} = \frac{\partial(F_1, \ldots, F_r)}{\partial(X_1, \ldots, X_r)}\) modulo \(m_A + (u, v, z, X)\). Hence, by the Implicit Function Theorem, there exists \(h := (h_1, \ldots, h_r) \in (m_A + (v, z))A(v, z)^r\) such that \(G(h) = 0\). Moreover \(f_1 = f_1 = h_1(v, \hat{z})\), thus we have \(f = g(v, \hat{z})\) where \(g(v, \hat{z}) := f_0 + h_1(v, \hat{z})\). This proves the lemma. \(\square\)
Then we can prove Theorem 3.5.8 by induction on $n$. First of all, since $A = B^T$ where $B$ is a complete regular local ring (by Cohen’s Structure Theorem), by using the same trick as in the proof of Corollary 3.2.7, we may replace $A$ by $B$ and assume that $A$ is a complete regular local ring. Let us assume that Theorem 3.5.8 is true for $n - 1$. We denote $x' := (x_1, ..., x_{n-1})$. We will denote by $y_1, ..., y_k$ the unknowns depending only on $x'$ and by $y_{k+1}, ..., y_m$ the unknowns depending on $x_n$. Let us consider the following system of equations

$$f(x', x_n, y_1(x'), ..., y_k(x'), y_{k+1}(x', x_n), y_m(x', x_n)) = 0.$$  

By Theorem 3.2.17 and Remark 3.2.22, we may assume that $\hat{y}_{k+1}, ..., \hat{y}_m \in k[[x']]$ and $\sum_{j \in \mathbb{N}} h_{i,j}(z)x_n^j \in k\langle z, x_n \rangle$ and $\hat{z} = (\hat{z}_1, ..., \hat{z}_s) \in (x')k[[]]^s$. We can write

$$f \left( x', x_n, y_1, ..., y_k, \sum_{j} h_{k+1,j}(z)x_n^j, ..., \sum_{j} h_{m,j}(z)x_n^j \right) = \sum_{j} G_j(x', y_1, ..., y_k, z)x_n^j$$

where $G_j(x', y_1, ..., y_k, z) \in k\langle x', y_1, ..., y_k, z \rangle$ for all $j \in \mathbb{N}$. Thus $\hat{y}_1, ..., \hat{y}_k, \hat{z}_1, ..., \hat{z}_s \in k[[x']]$ is a solution of the equations $G_j = 0$ for all $j \in \mathbb{N}$. Since $k(t, y_1, ..., y_k, z)$ is Noetherian, this system of equations is equivalent to a finite system $G_j = 0$ with $j \in E$ where $E$ is a finite subset of $\mathbb{N}$. Thus by the induction hypothesis applied to the system $G_j(x', y_1, ..., y_k, z) = 0$, $j \in E$, there exist $\tilde{y}_1, ..., \tilde{y}_k, \tilde{z}_1, ..., \tilde{z}_s \in k(x')$, with nested conditions, such that $\tilde{y}_i - \hat{y}_i, \tilde{z}_i \in (x')^s$, for $1 \leq i \leq k$ and $1 \leq l \leq s$, and $G_j(x', \tilde{y}_1, ..., \tilde{y}_k, \tilde{z}) = 0$ for all $j \in E$, thus $G_j(x', \tilde{y}_1, ..., \tilde{y}_k, \tilde{z}) = 0$ for all $j \in \mathbb{N}$.

Set $\hat{y}_i = \sum_{j \in \mathbb{N}} h_{i,j}(z)x_n^j$ for $k < j \leq m$. Then $\hat{y}_1, ..., \hat{y}_m$ satisfy the conclusion of the theorem.

$\square$

**Proposition 3.5.10.** — [ArMa65][AMR92] Let $A$ be a complete local normal domain and $v := (v_1, ..., v_n)$. If $f \in (m_A + (v))A(v)$ then there exists an integer $r \in \mathbb{N}$ and $F_r \in A[v][X_1, ..., X_r]$ such that $\frac{\partial(F_1, ..., F_r)}{\partial(X_1, ..., X_r)}$ is non-zero modulo $m_A + (v, X)$ and such that the unique $(f_1, ..., f_r) \in (m_A + (v))A(v)^r$ with $F(f_1, ..., f_r) = 0$ (according to the Implicit Function Theorem) is such that $f = f_1$.

**Proof.** — Let $P(v, X_1) \in A[v][X_1]$ be an irreducible polynomial such that $P(v, f) = 0$. Set $R := A[v][X_1]/(P(v, X_1))$ and let $\overline{R}$ be its normalization. Let $\varphi : R \rightarrow A(v)$ be the $A[v]$-morphism defined by $\varphi(X_1) = f$. Since $A(v)$ is an integrally closed domain, by the universal property of the normalization, the morphism $\varphi$ factors through $R \rightarrow \overline{R}$. Let $\overline{\varphi} : \overline{R} \rightarrow A(v)$ be the extension of $\varphi$ to $\overline{R}$. 


Since \( R \) is finitely generated over a local complete domain \( A \), then \( \overline{R} \) is module-finite over \( R \). Hence \( \overline{R} = \frac{A[x_1,x_2,\ldots,x_r]}{(F_1,\ldots,F_s)} \). Set \( f_i := \overline{\varphi}(X_i) \), for \( 2 \leq i \leq r \). By replacing \( X_i \) by \( X_i + a_i \) for some \( a_i \in A \) we may assume that \( f_i \in \mathfrak{m}_A + (v) \). Let us denote \( B := \overline{R}_{m_A+(v,x_1,\ldots,x_r)} \). Thus \( \overline{\varphi} \) induces a surjective \( A[v] \)-morphism \( B \to A[v] \) and by the universal property of the Henselization it induces a surjective \( A[v] \)-morphism \( B^h \to A(v) \) where \( B^h \) denotes the Henselization of \( B \). Moreover \( A[v]_{m_A+(v)} \to B \) induces a morphism between \( A(v) \) and \( B^h \) which is finite since \( A[v] \to \overline{R} \) is finite. Since \( B \) is an integrally closed local domain, then its completion is a local domain \([\text{Za48}]\), hence \( B^h \) is a local domain. If \( b \in B^h \) is in the kernel of \( B^h \to A(v) \), then \( b \) would satisfy \( b^k = 0 \) for some positive integer \( k \). But \( B^h \) being a domain, then \( b \) has to be zero. Thus \( B^h \to A(v) \) is injective hence \( B^h \) and \( A(v) \) are isomorphic. Moreover we have \( B^h \cong B \otimes_{A[v]_{m_A+(v)}} A(v) \). Using the definition of an étale morphism, since \( A[v]_{m_A+(v)} \to A(v) \) is faithfully flat, it is an exercise to check that \( A[v]_{m_A+(v)} \to B \) is étale. Thus \( s = r \) and \( \partial(F_1,\ldots,F_s) \) is non-zero modulo \( m_A + (v,X) \) and the unique solution of \( F = 0 \) in \((m_A + (v))A(v)^r \) is \( (f_1, f_2, \ldots, f_r) \).

Using ultraproducts methods we can deduce the following Strong Approximation result:

**Corollary 3.5.11.** — \([\text{BDLvdD79}]\) Let \( k \) be a field and \( f(x,y) \in k(x,y)^r \). There exists \( \beta : \mathbb{N} \to \mathbb{N} \) satisfying the following:

Let \( c \in \mathbb{N} \) and \( \overline{\varphi}(x) \in ((k[x])^m) \) satisfy \( f(x, \overline{\varphi}(x)) \in (x)^{\beta(c)} \). Let us assume that \( \overline{\varphi}_i(x) \in k[x_1,\ldots,x_{s_i}], \) \( 1 \leq i \leq m, \) for integers \( s_i, \) \( 1 \leq s_i \leq n \).

Then there exists a solution \( \tilde{\varphi}(x) \in ((k[x])^m) \) such that \( \overline{\varphi}_i(x) \in k[x_1,\ldots,x_{s_i}] \) and \( \overline{\varphi}(x) - \tilde{\varphi}(x) \in (x)^{\beta(c)} \).

**3.5.3. Nested Approximation in the analytic case.** — In the analytic case, Theorem 3.5.8 is no more valid, as shown in the following example:

**Example 3.5.12 (Gabrielov’s Example).** — \([\text{Gab71}]\) Let \( \varphi : \mathbb{C}\{x_1,x_2,x_3\} \to \mathbb{C}\{y_1,y_2\} \) be the morphism of analytic \( \mathbb{C} \)-algebras defined by \( \varphi(x_1) = y_1 \), \( \varphi(x_2) = y_1y_2 \), \( \varphi(x_3) = y_1e^{y_2} \).

Let \( f \in \text{Ker}(\overline{\varphi}) \) be written as \( f = \sum_{d=0}^{+\infty} f_d \) where \( f_d \) is a homogeneous polynomial of degree \( d \) for all \( d \in \mathbb{N} \). Then \( 0 = \overline{\varphi}(f) = \sum_{d} y_1^d f_d(1, y_2, y_2e^{y_2}) \).

Thus \( f_d = 0 \) for all \( d \in \mathbb{N} \) since \( 1, y_2, y_2e^{y_2} \) are algebraically independent over \( \mathbb{C} \). Hence \( \text{Ker}(\overline{\varphi}) = (0) \) and \( \text{Ker}(\varphi) = (0) \). This example is due to W. S. Osgood \([\text{Os16}]\).

- We may remark that \( "\varphi\left(x_3 - x_2e^{y_2}\right) = 0" \). But \( x_3 - x_2e^{y_2} \) is not an
element of \( \mathbb{C}\{x_1, x_2, x_3\} \).

Let us denote

\[
  f_n := \left( x_3 - x_2 \sum_{i=0}^{n} \frac{1}{i!} x_1^i \right) x_1^n \in \mathbb{C}[x_1, x_2, x_3], \forall n \in \mathbb{N}.
\]

Then

\[
  \varphi(f_n) = y_1^{n+1} y_2 \sum_{i=n+1}^{+\infty} \frac{y_2^i}{i!}, \forall n \in \mathbb{N}.
\]

Then we see that \((n + 1)! \varphi(f_n)\) is a convergent power series whose coefficients have module less than 1. Moreover if the coefficient of \(y^\text{expansion of } \varphi(\hat{g})\) is non zero then \(k = n + 1\). Thus \(h := \sum_n (n + 1)! \varphi(f_n)\) is a convergent power series since each of its coefficients has module less than 1. But \(\hat{g}\) being injective, the unique element whose image is \(h\) is necessarily

\[
  \hat{g} := \sum_n (n + 1)! f_n.
\]

and \(\sum_n (n + 1)! x_1^n\) is a divergent power series and \(\varphi(\hat{g}(x)) = h(y) \in \mathbb{C}\{y\}\). Hence \(\varphi(\mathbb{C}\{x\}) \subseteq \varphi(\mathbb{C}\{x\}) \cap \mathbb{C}\{y\}\).

\[\text{• By Lemma 3.5.1, } \hat{g}(x) = h(y) \text{ is equivalent to say that there exist } \hat{k}_1(x, y), \hat{k}_2(x, y), \hat{k}_3(x, y) \in \mathbb{C}\{x, y\} \text{ such that}\]

(3) \(\hat{g}(x) + (x_1 - y_1)\hat{k}_1(x, y) + (x_2 - y_1 y_2)\hat{k}_2(x, y) + (x_3 - y_1 e^{y_2})\hat{k}_3(x, y) - h(y) = 0.\)

Since \(\hat{g}(x)\) is the unique element whose image under \(\hat{g}\) equals \(h(y)\), Equation (3) has no convergent solution \(g(x) \in \mathbb{C}\{x, y\}, k_1(x, y), k_2(x, y), k_3(x, y) \in \mathbb{C}\{x, y\}\). Thus Theorem 3.5.8 is not true in the analytic setting.

Let us denote \(\hat{g}_1(x_1, x_2) := \sum_n (n + 1)! x_1^n\) and \(\hat{g}_2(x_1, x_2) := \hat{f}(x_1, x_2)\). By replacing \(y_1\) by \(x_1\), \(y_2\) by \(x_2\) and \(x_3\) by \(x_1 e^y\) in Equation (3) we see that the equation

(4) \(\hat{g}_1(x_1, x_2)x_1 e^y + \hat{g}_2(x_1, x_2) + (x_2 - x_1 y)\hat{k}(x, y) - h(x_1, y) = 0.\)

has a nested formal solution but no nested convergent solution.

Nevertheless there are, at least, three positive results about the nested approximation problem in the analytic category. They are the followings.
3.5. APPROXIMATION WITH CONSTRAINTS

3.5.3.1. Grauert’s Theorem. — The first one is due to H. Grauert who proved it in order to construct analytic deformations of a complex analytic germ in the case it has an isolated singularity. The approximation result of H. Grauert may be reformulated as: "if a system of complex analytic equations, considered as a formal nested system, admits an Artin function (as in Problem 2) which is the Identity function, then it has nested analytic solutions". We present here the result.

Set \( x := (x_1, \ldots, x_n), t := (t_1, \ldots, t_i), y = (y_1, \ldots, y_m) \) and \( z := (z_1, \ldots, z_l) \). Let \( f := (f_1, \ldots, f_r) \) be in \( \mathbb{C} \{ t, x, y, z \} \). Let \( I \) be an ideal of \( \mathbb{C} \{ t \} \).

**Theorem 3.5.13.** — Let \( d_0 \in \mathbb{N} \) and \( (\overline{y}(t), \overline{z}(t, x)) \in \mathbb{C}[t]^m \times \mathbb{C} \{ x \} \mathbb{C} \{ t \}^p \) satisfy
\[
\overline{y}(t, x, \overline{y}(t), \overline{z}(t, x)) \in I + (t)^{d_0}.
\]

Let us assume that for any \( d \geq d_0 \) and for any \( (y^{(d)}(t), z^{(d)}(t, x)) \in \mathbb{K}[t]^m \times \mathbb{K} \{ x \} \mathbb{K} \{ t \}^p \) such that, \( \overline{y}(t) - y^{(d)}(t) \in (t)^{d_0} \) et \( \overline{z}(t, x) - z^{(d)}(t, x) \in (t)^{d_0} \), and such that
\[
f(t, x, y^{(d)}(t), z^{(d)}(t, x)) \in I + (t)^d,
\]
there exists \( (\varepsilon(t), \eta(t, x)) \in \mathbb{K}[t]^m \times \mathbb{K} \{ x \} \mathbb{K} \{ t \}^p \) homogeneous in \( t \) of degree \( d \) such that
\[
f(t, x, y^{(d)}(t) + \varepsilon(t), z^{(d)}(t, x) + \eta(t, x)) \in I + (t)^{d+1}.
\]

Then there exists \( (\overline{y}(t), \overline{z}(t, x)) \mathbb{C} \{ t \}^m \times \mathbb{C} \{ x \} \mathbb{C} \{ t \}^p \) such that
\[
f(t, x, \overline{y}(t), \overline{z}(t, x)) \in I \quad \text{and} \quad \overline{y}(t) - \overline{y}(t, x), \overline{z}(t, x) - \overline{z}(t, x) \in (t)^{d_0}.
\]

The main ingredient of the proof is a result of Functional Analysis called "voisinages privilégiés" and proven by H. Cartan ([Ca44 Théorème α]). We do not give details here, but the reader may consult [JoPf00]. B. Malgrange generalized this result to partial differential equations in [Mal72].

3.5.3.2. Gabrielov’s Theorem. — The second positive result about the nested approximation problem in the analytic category is due to A. Gabrielov. Before giving his result, let us explain the context.

Let \( \varphi : A \rightarrow B \) be a morphism of analytic algebras where \( A := \mathbb{C} \{ x_1, \ldots, x_n \} \) and \( B := \mathbb{C} \{ y_1, \ldots, y_m \} \) are analytic algebras. Let us denote \( \varphi_i := \varphi(x_i) \) for \( 1 \leq i \leq n \). Let us denote by \( \widehat{\varphi} : \hat{A} \rightarrow \hat{B} \) the morphism induced by \( \varphi \). A. Grothendieck ([Gro60] and S. S. Abhyankar ([Ar71] raised the following question: Does \( \text{Ker}(\widehat{\varphi}) = \text{Ker}(\varphi) \))? Without loss of generality, we may assume that \( A = \mathbb{C} \{ x_1, \ldots, x_n \} \) and \( B = \mathbb{C} \{ y_1, \ldots, y_m \} \).

In this case, an element of \( \text{Ker}(\varphi) \) (resp. of \( \text{Ker}(\widehat{\varphi}) \)) is called an analytic (resp. formal) relation between \( \varphi_1(y), \ldots, \varphi_m(y) \). Hence the previous question is equivalent to the following: is any formal relation \( \tilde{S} \) between \( \varphi_1(y), \ldots, \varphi_n(y) \)
a linear combination of analytic relations?
This question is also equivalent to the following: may any formal relation between \(\varphi_1(y), \ldots, \varphi_n(y)\) be approximated by analytic relations for the \((x)\)-adic topology? In this form the problem is the "dual" problem to the Artin Approximation Problem.

In fact this problem is also a nested approximation problem. Indeed let \(\hat{S}\) be a formal relation between \(\varphi_1(y), \ldots, \varphi_n(y)\). This means that \(\hat{S}(\varphi_1(y), \ldots, \varphi_n(y)) = 0\).

By Lemma 3.5.1 this is equivalent to the existence of \(\hat{h}_1(x, y), \ldots, \hat{h}_n(x, y) \in \mathbb{C}[x, y]\) such that

\[
\hat{S}(x_1, \ldots, x_n) - \sum_{i=1}^{n} (x_i - \varphi_i(y))\hat{h}_i(x, y) = 0.
\]

If this equation has an analytic nested solution \(S(x) \in \mathbb{C}\{x\}, \ h_1(x, y), \ldots, h_n(x, y) \in \mathbb{C}\{x, y\}\), it gives an analytic relation between \(\varphi_1(y), \ldots, \varphi_n(y)\).

Example 3.5.14 — \([\text{Gab71}]\) Let us consider now the morphism

\[\psi : \mathbb{C}\{x_1, x_2, x_3, x_4\} \longrightarrow \mathbb{C}\{y_1, y_2\}\]

defined by

\[\psi(x_1) = y_1, \ \psi(x_2) = y_1 y_2, \ \psi(x_3) = y_1 y_2 e^{y_2}, \ \psi(x_4) = h(y_1, y_2)\]

Then \(x_4 - \hat{g}(x_1, x_2, x_3) \in \text{Ker}(\hat{\psi})\). On the other hand the morphism induced by \(\psi\) on \(\mathbb{C}[x_1, \ldots, x_4]/(x_4 - \hat{g}(x_1, x_2, x_3))\) is isomorphic to \(\hat{\psi}\) (where \(\varphi\) is the morphism of Example 3.5.12) that is injective. Thus we have \(\text{Ker}(\psi) = (x_4 - \hat{g}(x_1, x_2, x_3))\).

Since \(\text{Ker}(\psi)\) is a prime ideal of \(\mathbb{C}\{x\}\), \(\text{Ker}(\psi)\mathbb{C}[x]\) is a prime ideal of \(\mathbb{C}[x]\) included in \(\text{Ker}(\hat{\psi})\) by Proposition 3.4.1. Let us assume that \(\text{Ker}(\psi) \neq (0)\), then \(\text{Ker}(\psi)\mathbb{C}[x] = \text{Ker}(\hat{\psi})\) since \(\text{ht}(\text{Ker}(\hat{\psi})) = 1\). Thus \(\text{Ker}(\psi)\) is generated by one convergent power series denoted by \(f \in \mathbb{C}\{x_1, \ldots, x_4\}\) (in unique factorization domains, prime ideals of height one are principal ideals). Since \(\text{Ker}(\hat{\psi}) = (x_4 - \hat{g}(x_1, x_2, x_3))\), there exists \(u(x) \in \mathbb{C}[x]\), \(u(0) \neq 0\), such that \(f = u(x)(x_4 - \hat{g}(x_1, x_2, x_3))\). By applying Weierstrass Preparation Theorem to \(f\) with respect to \(x_4\) we see that \(u(x)\) and \(x_4 - \hat{g}(x_1, x_2, x_3)\) must be convergent, which is impossible since \(\hat{g}\) is a divergent power series. Hence \(\text{Ker}(\psi) = (0)\) but \(\text{Ker}(\hat{\psi}) \neq (0)\).

Nevertheless A. Gabrielov proved the following theorem:

Theorem 3.5.15 — \([\text{Gab73}]\) Let \(\varphi : A \longrightarrow B\) be a morphism of complex analytic algebras. Let us assume that the generic rank of the Jacobian matrix is equal to \(\text{dim}(A/\text{Ker}(\varphi))\).

Then \(\text{Ker}(\hat{\varphi}) = \text{Ker}(\varphi).\hat{A}\).
Sketch of the proof. — We give a sketch of the proof given by J.-Cl. Tougeron [To90]. As before we may assume that $A = \mathbb{C}\{x_1, ..., x_n\}$ and $B = \mathbb{C}\{y_1, ..., y_m\}$. Let us assume that $\ker(\varphi).\hat{A} \subset \ker(\hat{\varphi})$. Using a Bertini type theorem we may assume that $n = 3$, $\varphi$ is injective and $\dim(\frac{\mathbb{C}[x]}{\ker(\hat{\varphi})}) = 2$ (in particular $\ker(\hat{\varphi})$ is a principal ideal). Moreover, in this case we may assume that $m = 2$. After a linear change of coordinates we may assume that $\ker(\hat{\varphi})$ is generated by an irreducible Weierstrass polynomial of degree $d$ in $x_3$. Using change of coordinates and quadratic transforms on $\mathbb{C}\{y_1, y_2\}$ and using changes of coordinates of $\mathbb{C}\{x\}$ involving only $x_1$ and $x_2$, we may assume that $\varphi_1 = y_1$ and $\varphi_2 = y_1y_2$. Let us denote $f(y) := \varphi_3(y)$. Then we have
\[
f(y)^d + \sum_{i=1}^{d} \hat{a}_i(y_1, y_1y_2)f(y)^{d-1} + \cdots + \hat{a}_d(y_1, y_1y_2) = 0
\]
for some $\hat{a}_i(x) \in \mathbb{C}[x_1, x_2]$, $1 \leq i \leq d$. Then we want to prove that the $\hat{a}_i$'s may be chosen convergent in order to get a contradiction. Let us denote
\[
P(Z) := Z^d + \sum_{i=1}^{d} \hat{a}_i(x_1, x_2)Z^{d-1} + \cdots + \hat{a}_d(x_1, x_2) \in \mathbb{C}[x][Z].
\]
Since $\ker(\hat{\varphi})$ is prime we may assume that $P(Z)$ is irreducible. J.-Cl. Tougeron studies the algebraic closure $\bar{\mathbb{K}}$ of the field $\mathbb{C}(x_1, x_2)$. Let consider the following valuation ring
\[
V := \left\{ \frac{f}{g} / f, g \in \mathbb{C}(x_1, x_2), g \neq 0, \ord(f) \geq \ord(g) \right\},
\]
let $\hat{V}$ be its completion and $\bar{\mathbb{K}}$ the fraction field of $\hat{V}$. J.-Cl. Tougeron proves that the algebraic extension $\mathbb{K} \to \bar{\mathbb{K}}$ splits into $\mathbb{K} \to \mathbb{K}_1 \to \bar{\mathbb{K}}$ where $\mathbb{K}_1$ is a subfield of the following field
\[
\mathbb{L} := \left\{ A \in \bar{\mathbb{K}} / \exists \delta, a_i \in \mathbb{K}[x] \text{ is homogeneous } \forall i, \right\}
\]
\[
\ord \left( \frac{a_i}{\delta^{m(i)}} \right) = i, \exists a, b \text{ such that } m(i) \leq ai + b \quad \forall i \quad \text{and} \quad A = \sum_{i=0}^{\infty} \frac{a_i}{\delta^{m(i)}} \right\}.
\]
Moreover the algebraic extension $\mathbb{K}_1 \to \bar{\mathbb{K}}$ is the extension of $\mathbb{K}_1$ generated by all the roots of polynomials of the form $Z^q + g_1(x)Z^{q-1} + \cdots + g_q$ where $g_i \in \mathbb{C}(x)$ are homogeneous rational fractions of degree $ei$, $1 \leq i \leq q, e \in \mathbb{Q}$. A root of such polynomial is called a homogeneous element of degree $e$. For example, square roots of $x_1$ or of $x_1 + x_2$ are homogeneous elements of degree 2. We have $\bar{\mathbb{K}} \cap \mathbb{L} = \mathbb{K}_1$. In the same way he proves that the algebraic closure $\bar{\mathbb{K}}^{an}$ of $\mathbb{K}_1^{an}$, the fraction field of $\mathbb{C}\{x_1, x_2\}$ can be factorized as $\mathbb{K}^{an} \to \mathbb{K}_1^{an} \to \bar{\mathbb{K}}^{an}$ with $\mathbb{K}_1^{an} \subset \mathbb{L}^{an}$.
where
\[ L^{an} := \left\{ A \in \bar{K} / \exists \delta, a_i \in k[x] \text{ is homogeneous} \forall i, \text{ ord} \left( \frac{a_i}{\delta^{m(i)}} \right) = i, A = \sum_{i=0}^{\infty} a_i \delta^{m(i)} \right\} \]

\[ \exists a, b \text{ such that } m(i) \leq ai + b \forall i \text{ and } \exists r > 0 \text{ such that } \sum_i ||a_i|| r^i < \infty \]

and \[ ||a(x)|| := \max_{|z| \leq 1} |a(z_1, z_2)| \] for a homogeneous polynomial \( a(x) \).

Clearly, \( \xi := f(x_1, \frac{x_2}{x_1}) \) is an element of \( \bar{K} \) since it is a root of \( P(Z) \). Moreover \( \xi \) may be written \( \xi = \sum_{i=1}^\infty \xi_i \gamma^i \) where \( \gamma \) is a homogenous element and \( \xi_i \in L^{an} \cap K \) for any \( i \), i.e. \( \xi \in L^{an}[\gamma] \). Thus the problem is to show that \( \xi_i \in \bar{K}_i^{an} \) for any \( i \), i.e. \( L^{an} \cap K = \bar{K}^{an}_i \).

Then the idea is to resolve, by a sequence of blowing-ups, the singularities of the discriminant locus of \( P(Z) \) which is a germ of plane curve. Let us call \( \pi \) this resolution map. Then the discriminant of \( \pi(P)(Z) \) is normal crossing and \( \pi(P)(Z) \) defines a germ of hypersurface along the exceptional divisor of \( \pi \), denoted by \( E \). Let \( p \) be a point of \( E \). At this point \( \pi(P)(Z) \) may factor as a product of polynomials and \( \xi \) is a root of one of these factors denoted by \( Q_1(Z) \) and this root is a germ of an analytic function at \( p \). Then the other roots of \( Q_1(Z) \) are also in \( L^{an}[\gamma'] \) according to the Abhyankar-Jung Theorem, for some homogeneous element \( \gamma' \). Thus the coefficients of \( Q_1(Z) \) are in \( L^{an} \) and are analytic at \( p \).

Then the idea is to use the special form of the elements of \( L^{an} \) to prove that the coefficients of \( Q_1(Z) \) may be extended as analytic functions along the exceptional divisor \( E \) (the main ingredient in this part is the Maximum Principle).

We can repeat the latter procedure in another point \( p' \): we take the roots of \( Q_1(Z) \) at \( p' \) and using Abhyankar-Jung Theorem we construct new roots of \( \pi(P)(Z) \) at \( p' \) and the coefficients of \( Q_2(Z) := \prod_i (Z - \sigma_i) \), where \( \sigma_i \) runs over all these roots, are in \( L^{an} \) and are analytic at \( p' \). Then we extend the coefficients of \( Q_2(Z) \) everywhere along \( E \). Since \( \pi(P)(Z) \) has exactly \( d \) roots, this process stops after a finite number of steps. The polynomial \( Q(Z) := \prod_i (Z - \sigma_k) \), where the \( \sigma_k \)'s are the roots of \( \pi(P)(Z) \) that we have constructed, is a polynomial whose coefficients are analytic everywhere and it divides \( \pi(P)(Z) \). Thus, by Grauert’s Direct Image Theorem, there exists \( R(Z) \in \mathbb{C}\{x\}[Z] \) such that \( \pi(R)(Z) = Q(Z) \). Thus \( R(Z) \) divides \( P(Z) \), but since \( P(Z) \) is irreducible, then \( P(Z) = R(Z) \in \mathbb{C}\{x\}[Z] \) and the result is proven. \( \square \)

3.5.3.3. One variable Nested Approximation. — In the example of A. Gabrielov we can remark that the nested part of the solutions depends on two variables \( x_1 \) and \( x_2 \). In the case they depend only on one variable the nested approximation property is true. This is the following theorem:
**Theorem 3.5.16.** — (cf. Theorem 5.1 [DeLi80]) Let \( k \) be a field and let \( k[[x]] \) be a \( W \)-system over \( k \). Let \( t \) be one variable, \( x = (x_1, \ldots, x_n) \), \( y = (y_1, \ldots, y_m, y_{m+1}, \ldots, y_{m+k}) \), \( f \in k[[t, x, y]] \). Let \( \hat{y}_1, \ldots, \hat{y}_m \in (t)k[[y]] \) and \( \hat{y}_{m+1}, \ldots, \hat{y}_{m+k} \in (t, x)k[[t, x]] \) satisfy \( f(t, x, \hat{y}) = 0 \). Then, for any \( c \in \mathbb{N} \), there exists \( \hat{y}_1, \ldots, \hat{y}_m \in (t)k[[y]] \), \( \hat{y}_{m+1}, \ldots, \hat{y}_{m+k} \in (t, x)k[[t, x]] \) such that \( f(t, x, \hat{y}) = 0 \) and \( \hat{y} - \hat{y} \in (t, x)^c \).

**Example 3.5.17.** — The main example is the case where \( k \) is a valued field and \( k[[x]] \) is the ring of convergent power series over \( k \).

**Proof.** — The proof is very similar to the proof of Theorem 3.5.8. Set \( u := (u_1, \ldots, u_j), j \in \mathbb{N} \) and Set

\[
k[[t]][(u)] := \{ f(z_1(t), \ldots, z_s(t), u) \in k[[t, u]] / f(z_1, \ldots, z_s, u) \in k[[z, u]] \text{ and } z_1(t), \ldots, z_s(t) \in (t)k[[u]] \text{ for some } s}.\]

The rings \( k[[t]][(u)] \) form a \( W \)-system over \( k[[t]] \) (cf. Lemma 52, [DeLi80] but it is straightforward to check it since \( k[[x]] \) is a \( W \)-system over \( k \) - in particular, if \( \text{char}(k) > 0 \), \( v_i \) of Definition 3.2.11 is satisfied since \( \nu \) of Definition 3.2.11 is satisfied for \( k[[x]] \)). By Theorem 3.2.14 applied to

\[
f(t, \hat{y}_1, \ldots, \hat{y}_m, y_{m+1}, \ldots, y_{m+k}) = 0\]

there exist \( \bar{y}_{m+1}, \ldots, \bar{y}_{m+k} \in k[[t]][(x)] \) such that \( f(t, \hat{y}_1, \ldots, \hat{y}_m, \bar{y}_{m+1}, \ldots, \bar{y}_{m+k}) = 0 \) and \( \bar{y}_n - \hat{y}_n \in (t, x)^c \) for \( m < i \leq m + k \).

Let us write \( \bar{y}_i = \sum_{\alpha \in \mathbb{N}^n} h_{i, \alpha}(\bar{z})x^\alpha \) with \( \sum_{\alpha \in \mathbb{N}^n} h_{i, \alpha}(z)x^\alpha \in k[[z, x]] \) and \( \bar{z} = (\bar{z}_1, \ldots, \bar{z}_s) \in k[[t]]. \) We can write

\[
f(t, x_1, \ldots, y_m, \sum_{\alpha} h_{m+1, \alpha}(z)x^\alpha, \ldots, \sum_{\alpha} h_{m+k, \alpha}(z)x^\alpha) = \sum_{\alpha} G_{\alpha}(t, y_1, \ldots, y_m, z)x^\alpha\]

where \( G_{\alpha}(t, y_1, \ldots, y_m, z) \in k[[t, y_1, \ldots, y_m, z]] \) for all \( \alpha \in \mathbb{N}^n \). Thus \( \hat{y}_1, \ldots, \hat{y}_m, \bar{z}_1, \ldots, \bar{z}_s \in k[[t]] \) is a solution of the equations \( G_{\alpha} = 0 \) for all \( \alpha \in \mathbb{N}^n \). Since \( k[[t, y_1, \ldots, y_m, z]] \) is Noetherian, this system of equations is equivalent to a finite system \( G_{\alpha} = 0 \) with \( \alpha \in E \) where \( E \) is a finite subset of \( \mathbb{N}^n \). Thus by Theorem 3.2.14 applied to the system \( G_{\alpha}(t, y_1, \ldots, y_m, z) = 0 \), \( \alpha \in E \), there exist \( \hat{y}_1, \ldots, \hat{y}_m, \bar{z}_1, \ldots, \bar{z}_s \in k[[t]] \) such that \( \bar{y}_i - \hat{y}_i, \bar{z}_j - \hat{z}_j \in (t)^c \) for \( 1 \leq i \leq m \) and \( 1 \leq j \leq s \), and \( G_{\alpha}(t, \hat{y}_1, \ldots, \hat{y}_m, \bar{z}) = 0 \) for all \( \alpha \in E \), thus \( G_{\alpha}(t, \hat{y}_1, \ldots, \hat{y}_m, \bar{z}) = 0 \) for all \( \alpha \in \mathbb{N}^n \).

Set \( \bar{y}_i = \sum_{\alpha \in \mathbb{N}^n} h_{i, \alpha}(\bar{z})x^\alpha \) for \( m < i \leq m + k \). Then \( \hat{y}_1, \ldots, \hat{y}_{m+k} \) satisfy the conclusion of the theorem.

**Remark 3.5.18.** — The proof of this theorem uses in an essential way the Weierstrass Division Property (in order to show that \( k[[t]][(u)] \) is a Noetherian local ring, which is the main condition to use Theorem 3.2.17) The Henselian
and excellent conditions may be proven quite easily). It is an open question to know if this result remains true if we do no have the Weierstrass Division Property.

For example let $C_n$ be the ring of germs of $k$-valued Denjoy-Carleman functions defined at the origin of $\mathbb{R}^n$, where $k = \mathbb{R}$ or $\mathbb{C}$ (see [Th08] for definitions and properties of these rings). It is still an open problem to know if $C_n$ is Noetherian or not for $n \geq 2$ ($C_1$ is a discrete valuation ring, thus it is Noetherian). These rings have similar properties to the Weierstrass systems (stability by partial derivatives, stability by division by coordinates, ...), except that there is no Weierstrass Division Theorem [Chi76]. For instance, there exists $f \in C_1$ and $\hat{g} \in \mathbb{k}[x] \setminus C_1$ such that $f(x) = \hat{g}(x^2)$ (see [Th08]). This implies that

$$f(x) = (x^2 - y)\hat{h}(x, y) + \hat{g}(y)$$

where $\hat{h}(x, y) \in \mathbb{k}[x, y]$ but Equation (5) has no nested solution in $C_1 \times C_2$.

On the other hand, if the rings $C_n$ were Noetherian, since their completions are regular local rings, they would be regular. Then using Example 3.7.4 iii) we see that they would be excellent (see also [ElKh12]). Thus these rings would satisfy Theorem 3.2.17. It would show that the Weierstrass Division Theorem is necessary to obtain Theorem 3.5.16.

3.5.4. Other examples of approximation with constraints. — We present here some examples of positive or negative answers to Problems 1 and 2 in several contexts.

Example 3.5.19. — [Mi78b] P. Milman proved the following theorem:

**Theorem 3.5.20.** Let $f \in \mathbb{C}\{x, y, u, v\}$ where $x := (x_1, ..., x_n)$, $y := (y_1, ..., y_n)$, $u := (u_1, ..., u_m)$, $v := (v_1, ..., v_m)$. Then the set of convergent solutions of the following system:

$$(6) \begin{cases} f(x, y, u(x, y), v(x, y)) = 0 \\ \frac{\partial u_k}{\partial x_j}(x, y) - \frac{\partial v_k}{\partial y_j}(x, y) = 0 \\ \frac{\partial v_k}{\partial x_j}(x, y) + \frac{\partial u_k}{\partial y_j}(x, y) = 0, \quad 1 \leq j \leq n \end{cases}$$

is dense (for the $(x, y)$-adic topology) in the set of formal solutions of this system.

**Hints on the proof.** Let $(\hat{u}(x, y), \hat{v}(x, y)) \in \mathbb{C}\{x, y\}^{2m}$ be a solution of (6). Let us denote $z := x + iy$ and $w := u + iv$. In this case the Cauchy-Riemann equations of (6) are equivalent to $\hat{w}(z, \overline{z}) := \hat{u}(x, y) + i\hat{v}(x, y) \in \mathbb{C}\{z\}$ (or in
3.5. APPROXIMATION WITH CONSTRAINTS

Let \( \varphi : \mathbb{C}\{z, z, w\} \rightarrow \mathbb{C}[z, z] \) and \( \psi : \mathbb{C}\{z, w\} \rightarrow \mathbb{C}[w] \) be the morphisms defined by

\[
\varphi(h(z, z, w, \overline{w}(z))) := h(z, z, \hat{w}(z), \overline{\hat{w}(z)}) \quad \text{and} \quad \psi(h(z, w)) := h(z, \hat{w}(z)).
\]

Milman proved that \( \ker(\varphi) = \ker(\psi) \mathbb{C}\{z, z, w, \overline{w}\} + \ker(\psi) \mathbb{C}\{z, z, w, \overline{w}\} \).

Since \( \ker(\psi) \) (as an ideal of \( \mathbb{C}\{z, w\} \)) satisfies Theorem 3.2.1, the result follows.

This proof does not give the existence of an Artin function for this kind of system, since the proof consists in reducing Theorem 3.5.20 to Theorem 3.2.1, but this reduction depends on the formal solution of (6). Nevertheless in [Hic-Ro11], it is proven that such a system admits an Artin function using ultraproducts methods. The survey [Mir13] is a good introduction for application of Artin Approximation in CR geometry.

**Example 3.5.21.** — [BiMi79] Let \( G \) be a reductive algebraic group. Suppose that \( G \) acts linearly on \( \mathbb{C}^n \) and \( \mathbb{C}^m \). We say that \( y(x) \in \mathbb{C}[x]^m \) is equivariant if \( y(\gamma x) = \gamma y(x) \) for all \( \gamma \in G \).

E. Bierstone and P. Milman proved that, in Theorem 3.2.1, the constraint for the solutions of being equivariant may be preserved for convergent solutions:

**Theorem 3.5.22.** — [BiMi79] Let \( f(x, y) \in \mathbb{C}\{x, y\}^r \). Then the set of equivariant convergent solutions of \( f = 0 \) is dense in the set of equivariant formal solutions of \( f = 0 \) for the \((x)\)-adic topology.

This result remains true if we replace \( \mathbb{C} \) (resp. \( \mathbb{C}\{x\} \) and \( \mathbb{C}\{x, y\} \)) by any field of characteristic zero \( k \) (resp. \( k(x) \) and \( k(x, y) \)).

Using ultraproducts methods we may probably prove that Problem 2 has a positive answer in this case.

**Example 3.5.23.** — [BDLvdD79] Let \( k \) be a field. Let us consider the following differential equation:

\[
(7) \quad a^2 x_1 \frac{\partial f}{\partial x_1}(x_1, x_2) - x_2 \frac{\partial f}{\partial x_2}(x_1, x_2) = \sum_{i,j \geq 1} x_1^i x_2^j \left( \frac{x_1}{1 - x_1} \right) \left( \frac{x_2}{1 - x_2} \right).
\]

For \( a \in \mathbb{k}, a \neq 0 \), this equation has only the following solutions

\[
f(x_1, x_2) := b + \sum_{i,j \geq 1} \frac{x_1^i x_2^j}{a^{i-j}}, \quad b \in \mathbb{k}.
\]
Let us consider the following system of equations:

\[
\begin{align*}
\sum_{i,j\geq 1} x_i^i x_j^j &= y_0^2 x_1 y_5 (x_1, x_2) - x_2 y_7 (x_1, x_2) \\
y_1 (x_1, x_2) &= y_2 (x_3, x_4, x_5) + (x_1 - x_3) z_1 (x) + (x_2 - x_4) z_2 (x) \\
y_2 (x_3, x_4, x_5) &= y_1 (x_1, x_2) + x_5 y_5 (x_1, x_2) + x_5^2 y_6 (x_3 - x_1 - x_5) z_3 (x) + (x_4 - x_2) z_4 (x) \\
y_3 (x_3, x_4, x_5) &= y_1 (x_1, x_2) + x_5 y_7 (x_1, x_2) + x_5^2 y_8 (x_3 - x_1) z_5 (x) + (x_4 - x_2 - x_5) z_5 (x) \\
y_0 (x_1, x_2) &= y_10 (x_3, 4, x_5) \quad \text{ i.e. } y_9 \in k \text{ and } y_9 y_{11} = 1.
\end{align*}
\]

It is straightforward, by Lemma 3.5.1 and Example 3.5.2, to check that \((a, f (x_1, x_2))\) is a solution of (7) if and only if (8) has a solution when modulo \(c\).

Moreover, if \(y_1, \ldots, y_{11}, z_1, \ldots, z_5\) is a solution of (8) with \(y_0 = 0\), then \((y_9, y_1)\) is a solution of (7).

Thus (8) has no solution in \(\mathbb{Q}[[x]]\). But clearly, (7) has solutions in \(\mathbb{Q}[[x]] / (x)^c\) for any \(c \in \mathbb{N}\) and the same is true for (8). This shows that Proposition 3.3.24 is not valid if the base field is not \(\mathbb{C}\).

**Example 3.5.24.** — BDLvdD79 Let us assume that \(k = \mathbb{C}\) and consider the latter example. The system of equations (8) does not admit an Artin function. Indeed, for any \(c \in \mathbb{N}\), there is \(a_c \in \mathbb{Q}\), such that (8) has a solution modulo \((x)^c\) with \(y_9 = a_c\). But there is no solution in \(\mathbb{C}[[x]]\) with \(y_9 = a_c\) modulo \((x)\), otherwise \(y_0 = a_c\) which is not possible.

Thus systems of equations with constraints does not satisfy Problem 2 in general.

### 3.6. Weierstrass Preparation Theorem

In this part set \(x := (x_1, \ldots, x_n)\) and \(x' := (x_1, \ldots, x_{n-1})\). Moreover \(k\) will denote a local ring of maximal ideal \(m\) (if \(k\) is a field, \(m = (0)\)). A local subring of \(k[[x]]\) will be a subring of \(k[[x]] / (x)^c\) which is a local ring and whose maximal ideal is generated by \((m + (x)) \cap A\).

**Definition 3.6.1.** — If \(f \in k[[x]]\) we say that \(f\) is regular of order \(d\) with respect to \(x_n\) if \(f = u x_n^d\) modulo \(m + (x')\) where \(u\) is invertible in \(k[[x]] / m + (x') \simeq \frac{k}{m}[[x_n]]\).

**Definition 3.6.2.** — Let \(A\) be a local subring of \(k[[x]]\). We say that \(A\) has the Weierstrass Division Property if for any \(f, g \in A\) such that \(f\) is regular of
order \(d\) with respect to \(x_n\), there exist \(q \in A\) and \(r \in (A \cap \k[[x']])[x_n]\) such that \(\deg_{x_n}(r) < d\) and \(g = qf + r\).

**Definition 3.6.3.** — Let \(A\) be a local subring of \(\k[[x]]\). We say that \(A\) satisfies the Weierstrass Preparation Theorem if for any \(f \in A\) which is regular with respect to \(x_n\), there exist an integer \(d\), a unit \(u \in A\) and \(a_1(x'), \ldots, a_d(x') \in A \cap (x')\k[[x]]\) such that

\[ f = u \left( x_n^d + a_1(x')x_n^{d-1} + \cdots + a_d(x') \right). \]

In this case \(f\) is necessarily regular of order \(d\) with respect to \(x_n\).

**Remark 3.6.4.** — Clearly, if they exist, \(q\) and \(r\) are unique in Definition 3.6.2. The same is true for \(u\) and the \(a_i(x')\)'s in Definition 3.6.3.

**Lemma 3.6.5.** — If a local subring \(A\) of \(\k[[x]]\) has the Weierstrass Division Property then it satisfies the Weierstrass Preparation Theorem.

**Proof.** — If \(A\) has the Weierstrass Division Property and if \(f \in A\) is regular of order \(d\) with respect to \(x_n\), then we can write \(x_n^d = qf + r\) where \(r \in (A \cap \k[[x']])[x_n]\) such that \(\deg_{x_n}(r) < d\). Thus \(qf = x_n^d - r\). Since \(f\) is regular of order \(d\) with respect to \(x_n\), then \(q\) is invertible in \(\k[[x]]\) and \(r \in (\mathfrak{m} + (x'))\). Thus \(q \notin (\mathfrak{m} + (x))\) and \(q\) is invertible in \(A\). Hence \(f = q^{-1}(x_n^d - r)\).

**Theorem 3.6.6.** — The following rings have the Weierstrass Division Property:

i) The ring \(A = \k[x]\) where \(\k\) is complete local ring (Bo65).

ii) The ring \(A = \k(x)\) of algebraic power series where \(\k\) is a field or a Noetherian Henselian local ring of characteristic zero which is analytically normal (Laf65 and Laf67).

iii) The ring \(A = \k\{x\}\) of convergent power series over a valued field \(\k\) (Na62).

## 3.7. Regular morphisms and excellent rings

**Definition 3.7.1.** — Let \(\varphi : A \to B\) be a morphism of Noetherian rings. We say that \(\varphi\) is regular if it is flat and if for any prime ideal \(P\) of \(A\), the \(\k(P)\)-algebra \(B \otimes_A \k(P)\) is geometrically regular (where \(\k(P) := \kappa(\mathfrak{p}_A)\) is the residue field of \(A\)). This means that \(B \otimes_A \k\) is a regular Noetherian ring for any finite field extension of \(\k(P)\).

**Example 3.7.2.** —
i) si $A$ et $B$ sont des corps, alors $A \rightarrow B$ est régulier si et seulement si $B$ est une extension séparable de $A$.

ii) Si $A$ est excellent, pour tout idéal $I$ de $A$, l’morphisme $A \rightarrow \hat{A}$ est régulier où $\hat{A} := \lim_{\leftarrow I} A$ est l’idéal-complétion de $A$ (cf. [GrDi65] 7.8.3).

iii) Si $V$ est une anneau d’évaluation discrète, alors la morphisme $V \rightarrow \hat{V}$ est régulier si et seulement si Frac$(V) \rightarrow$ Frac$(\hat{V})$ est séparable. En effet, $V \rightarrow \hat{V}$ est toujours plate et cette morphisme induit un isomorphisme sur les corps résiduels.

iv) Soit $X$ un compact manifolds de Nash, soit $N(X)$ le corps des fonctions de Nash sur $X$ et soit $O(X)$ le corps des fonctions analytiques sur $X$. Alors l’inclusion naturelle $N(X) \rightarrow O(X)$ est régulier (cf. [CRS95]).

v) Soit $L \subset \mathbb{C}^n$ un compact polyédre polynomial et $B$ le corps des fonctions holomorphes à $L$. Alors la morphisme des constantes $C \rightarrow B$ est régulier (cf. [Le95]). Ce résultat et le précédent permettent d’utiliser le théorème 3.2.16 pour obtenir des résultats d’approximation globale en géométrie complexe ou réelle.

En ce qui concerne l’approximation d’Artin, nous serons principalement intéressés par la morphisme $A \rightarrow \hat{A}$. Nous devons savoir ce qu’est un anneau excellent.

**Définition 3.7.3.** — Un anneau excellent $A$ est excellent si les conditions suivantes sont remplies:

i) $A$ est catégoriquement universel.

ii) Pour tout $p \in \text{Spec}(A)$, l’orbite formelle de $A_p$ est géométriquement régulier.

iii) Pour tout $p \in \text{Spec}(A)$ et pour tout extension séparable finie Frac$(A_p) \rightarrow \mathbb{K}$, il existe un anneau $B$ finiment engendré d’une $A_p$-algèbre, contenant $A_p$, et un tel que Frac$(B) = \mathbb{K}$ et le ensemble des points réguliers de Spec$(B)$ contient un ensemble ouvert non-void.

Cette définition peut être un peu obscure d’abord vue. Nous donnons ici les exemples principaux d’anneaux excellents:

**Exemple 3.7.4.** —

i) Tous les anneaux locaux complets (à savoir n’importe quel corps) sont excellents. Les anneaux de Dedekind de caractéristique nulle sont excellents. Toute anneau qui est essentiellement de type fini sur un excellent anneau est excellent. ([GrDi65] 7-8-3).

ii) Si $k$ est un corps valué complet, alors $k\{x_1, ..., x_n\}$ est excellent [Ki69].

iii) Nous avons le suivant résultat: Soit $A$ une anneau régulier contenant un corps de caractéristique nulle dénoté par $k$. Supposons que pour un idéal maximal $m$, le corps extérieur $k \rightarrow \frac{A}{m}$ est algébrique et ht$(m) = n$. Supposons outre que $D_1, ..., D_n \in \text{Der}_k(A)$ et $x_1, ..., x_n \in A$ tel que $D_i(x_j) = \delta_{i,j}$. Alors $A$ est excellent (cf. Théorème 102 [Mat80]).
iv) A Noetherian local ring $A$ is excellent if and only if it is universally catenary and $A \to \hat{A}$ is regular ([GrDi65 7-8-3 i]). In particular, if $A$ is a quotient of a local regular ring, then $A$ is excellent if and only if $A \to \hat{A}$ is regular (cf. [GrDi65 5-6-4]).

**Example 3.7.5.** — [Na62, Mat80] Let $k$ be a field of characteristic $p > 0$ such that $[k : k^p] = \infty$ (for instance let us take $k = \mathbb{F}_p(t_1, \ldots, t_n, \ldots)$). Let $V := k^p[[x]][k]$ where $x$ is a single variable, i.e. $V$ is the ring of power series $\sum_{i=0}^{\infty} a_i x^i$ such that $[k^p(a_0, a_1, \ldots) : k^p] < \infty$. Then $V$ is a discrete valuation ring whose completion is $k[[x]]$ and it is a Henselian ring. We have $\hat{V}^p \subset V$, thus $\text{Frac}(\hat{V}) : \text{Frac}(V)$ is purely inseparable. Hence $V \to \hat{V}$ is not regular by Example 3.7.2 and $V$ is not excellent by Example 3.7.4 iv).

On the other hand, let $f$ be the power series $\sum_{i=0}^{\infty} a_i x^i$, $a_i \in k$ such that $[k^p(a_0, a_1, \ldots) : k^p] = \infty$. Then $f \in \hat{V}$ but $f \notin V$, and $f^p \in V$. Thus $f$ is the only root of the polynomial $y^p - f^p$. This shows that the polynomial $y^p - f^p \in V[y]$ does not satisfies Theorem 3.2.16.

### 3.8. Étale morphisms and Henselian rings

The material presented here is very classical and has first been studied by G. Azumaya and M. Nagata. We will give a quick review of the definitions and properties that we need for the understanding of the rest of the chapter. Nevertheless, the reader may consult [Na62, GrDi65, Ra70] or [Iv73].

**Example 3.8.1.** — In classical algebraic geometry, the Zariski topology has too few open sets. For instance, there is no Implicit Function Theorem. Let $X$ be the zero set of the polynomial $y^2 - x^2(x+1)$ in $\mathbb{C}^2$. On an affine open neighborhood of 0, denoted by $U$, $X \cap U$ is equal to $X$ minus a finite number of points, thus $X \cap U$ is irreducible since $X$ is irreducible. In the analytic topology, we can find an open neighborhood of 0, denoted by $U$, such that $X \cap U$ is reducible, for instance take $U = \{(x, y) \in \mathbb{C}^2 / |x|^2 + |y|^2 < 1/2\}$. This comes from the fact that $x^2(1+x)$ is the square of an analytic function defined on $U \cap (\mathbb{C} \times \{0\})$. Let $z(x)$ be such an analytic function, $z(x)^2 = x^2(1+x)$.

In fact we can obtain $z(x)$ from the Implicit Function Theorem. We see that $z(x)$ is a root of the polynomial $Q(x, z) := z^2 - x^2(1+x)$. We have $Q(0, 0) = 0$, thus we can use directly the Implicit Function Theorem to obtain $z(x)$ from its minimal polynomial. Nevertheless let us take $P(x, t) := (t+1)^2 - (1+x) = t^2 + 2t - x$. Then $P(0, 0) = 0$ and $\frac{\partial P}{\partial t}(0, 0) = 2 \neq 0$. Thus, from the Implicit function Theorem, there exists $t(x)$ analytic on a neighborhood of 0 such that $t(0) = 0$ and $P(x, t(x)) = 0$. If we denote $z(x) := x(1+t(x))$, we have $z^2(x) = x^2(1+x)$. In
fact \( z(x) \in B := \mathbb{C}[x,t]_{(x,t)}^{(P(x,t))} \). The morphism \( \mathbb{C}[x] \rightarrow B \) is an example of étale morphism.

**Definition 3.8.2.** — Let \( \varphi : A \rightarrow B \) be a ring morphism. We say that \( \varphi \) is smooth (resp. étale) if for any \( A \)-algebra \( C \) along with an ideal \( I \) such that \( I^2 = (0) \) and any morphism of \( A \)-algebras \( \psi : B \rightarrow \mathbb{C}[t]_T \) there exists a morphism \( \sigma : B \rightarrow C \) (resp. a unique morphism) such that the following diagram commutes:

\[
\begin{array}{ccc}
a & \varphi & \rightarrow & b \\
\downarrow & \downarrow & \sigma & \downarrow \\
c & \rightarrow & c \\
\end{array}
\]

**Example 3.8.3.** — Let \( k := \mathbb{R} \) or \( \mathbb{C} \) and let us assume that \( A = k[x_1,\ldots,x_n]_J \) and \( B = A[y_1,\ldots,y_m]_K \) for some ideals \( J \) and \( K \). Let \( X \) be the zero locus of \( J \) in \( k^n \) and \( Y \) be the zero locus of \( K \) in \( k^n + m \). The morphism \( \varphi : A \rightarrow B \) defines a regular map \( \Phi : Y \rightarrow X \). Let \( C := \mathbb{C}[t]/(t^2) \) and \( I := (t) \). Let \( f_1(x),\ldots,f_r(x) \) be generators of \( J \).

A morphism \( A \rightarrow C \) is given by elements \( a_i, b_i \in k \) such that \( f_j(a_1 + b_1t,\ldots,a_n + b_nt) \in (t)^2 \) for \( 1 \leq j \leq r \). We have

\[
f_j(a_1 + b_1t,\ldots,a_n + b_nt) = f_j(a_1,\ldots,a_n) + \left( \sum_{i=1}^n \frac{\partial f_j}{\partial x_i}(a_1,\ldots,a_n)b_i \right) t \text{ mod. } (t)^2.
\]

Thus a morphism \( A \rightarrow C \) is given by a point \( x := (a_1,\ldots,a_n) \in X \) (i.e. such that \( f_j(a_1,\ldots,a_n) = 0 \) for all \( j \)) and a tangent vector \( u := (b_1,\ldots,b_n) \) to \( X \) at \( x \) (i.e. such that \( \sum_{i=1}^n \frac{\partial f_j}{\partial x_i}(a_1,\ldots,a_n)b_i = 0 \) for all \( j \)). In the same way a \( A \)-morphism \( B \rightarrow \mathbb{C}[t]/(t^2) \) is given by a point \( y \in Y \). Moreover the first diagram commutes if and only if \( \Phi(y) = x \).

Then \( \varphi \) is smooth if for any \( x \in X \), any \( y \in Y \) and any tangent vector \( u \) to \( X \) at \( x \) such that \( \Phi(y) = x \), there exists a tangent vector \( v \) to \( Y \) at \( y \) such that \( D_y(\Phi)(v) = u \). And \( \varphi \) is étale if and only if \( v \) is unique. This shows that smooth morphisms correspond to submersions and étale morphisms to local diffeomorphisms.

**Example 3.8.4.** — Let \( \varphi : A \rightarrow B_p \) be the canonical morphism where \( B := \mathbb{C}[x]_{(P(x))} \) and \( p \) is a prime ideal of \( B \) such that \( \frac{\partial P}{\partial x}(x) \notin p \). If we have such
3.8. Étale Morphisms and Henselian Rings

A commutative diagram

\[
\begin{array}{ccc}
A & \xrightarrow{\varphi} & B_p \\
\downarrow & & \downarrow \psi \\
C & \xrightarrow{I} & C_T
\end{array}
\]

then the morphism \( B_p \to C_T \) is given by an element \( c \in C \) such that \( P(c) \in I \).

Looking for a lifting of \( \psi \) is equivalent to finding \( \varepsilon \in I \) such that \( P(c + \varepsilon) = 0 \).

We have

\[
P(c + \varepsilon) = P(c) + \frac{\partial P}{\partial x}(c)\varepsilon
\]

since \( I^2 = (0) \). Since \( \frac{\partial P}{\partial x} \) is invertible in \( B_p \), \( \frac{\partial P}{\partial x}(c) \) is invertible in \( C_T \), i.e. there exists \( a \in C \) such that \( a\frac{\partial P}{\partial x}(c) = 1 \mod. I \). Moreover \( a \) is unique modulo \( I \).

For any \( \eta \in I \) let us set \( \varepsilon := -P(c)(a + \eta) \). Since \( P(c) \in I \), \( \varepsilon \) does not depend on \( \eta \) and the lifting of \( \psi \) is unique. This proves that \( \varphi \) is étale. Compare this example with Example 3.8.1.

Definition 3.8.5. — Étale morphisms of Example 3.8.4 are called standard étale morphisms. We can prove that if \( A \) and \( B \) are local rings then any étale morphism is standard ([IV73] III. 2).

Example 3.8.6 (Jacobian Criterion). — We can generalize the former example as follows. If \( k \) is a field and \( \varphi : k \to B := \k[x_1, ..., x_n]_{(g_1, ..., g_r)} \) where \( m := (x_1 - c_1, ..., x_n - c_n) \) then \( \varphi \) is smooth if and only if the Jacobian matrix \( \left( \frac{\partial g_i}{\partial x_j}(c) \right) \) has rank equal to the height of \( (g_1, ..., g_r) \). This is equivalent to say that \( V(I) \) has a non-singular point at the origin. Let us recall that the fibers of submersions are always smooth.

Definition 3.8.7. — Let \( A \) be a local ring. An étale neighbourhood of \( A \) is an étale local morphism \( A \to B \) inducng an isomorphism on the residue fields.

If \( A \) is a local ring, the étale neighbourhoods of \( A \) form a filtered inductive limit and the limit of this system is called the Henselization of \( A \) ([IV73] III. 6. or [Ra69] VIII) and denoted by \( A^h \).

We say that \( A \) is Henselian if \( A = A^h \). The morphism \( A \to A^h \) is universal among all the morphisms \( A \to B \) inducing an isomorphisms on the residue fields and where \( B \) is Henselian.

Proposition 3.8.8. — If \( A \) is a Noetherian local ring, then its Henselization \( A^h \) is a Noetherian local ring and \( A \to A^h \) is faithfully flat. If \( \varphi : A^h \to B \) is an étale neighbourhood of \( A^h \), then there is a section \( \sigma : B \to A \), i.e. \( \sigma \circ \varphi = \text{id}_{A^h} \).
Example 3.8.15 of convergent power series.

universal property of the Henselization.

functions at the origin of is a Henselian local ring but it is not Noetherian. The ring of germ of analytic then there exists \( s = s \mod m^c \).

Proof. — Since \( A \) is Henselian and \( \varphi \) is étale then \( A \) is isomorphic to the Henselization of \( B \). Moreover \( \frac{\Delta}{\Delta_A} \) is Henselian. The result comes from the universal property of the Henselization.

Definition 3.8.10. — Let \( A \) be a Henselian local ring and \( x := (x_1, ..., x_n) \). Then the Henselization of \( A[x]\_{m_A+(x)} \) is denoted by \( A(x) \).

Remark 3.8.11. — Let \( P(y) \in A[y] \) and \( a \in A \) satisfy \( P(a) \in m_A \) and \( \frac{\partial P}{\partial y}(a) \notin m_A \). If \( A \) is Henselian, then \( A \approx \frac{A[y]}{(P(y))} m_A+(y-a) \) is an étale neighborhood of \( A \), thus it admits a section. This means that there exists \( \tilde{y} \in m_A \) such that \( P(a + \tilde{y}) = 0 \).

If \( A \) is a local ring, then any étale neighborhood of \( A \) is of the previous form. Thus, by Proposition 3.8.8 we have the following proposition:

Proposition 3.8.12. — Let \( A \) be a local ring. Then \( A \) is Henselian if and only if for any \( P(y) \in A[y] \) and \( a \in A \) such that \( P(a) \in m_A \) and \( \frac{\partial P}{\partial y}(a) \notin m_A \) there exists \( \tilde{y} \in m_A \) such that \( P(a + \tilde{y}) = 0 \).

We can generalize this proposition as follows:

Theorem 3.8.13 (Implicit Function Theorem). — Let \( f(y) \in A[y]^r \), \( y = (y_1, ..., y_m) \), \( r \leq m \). Let \( J \) be the ideal of \( A[y] \) generated by the \( r \times r \) minors of the Jacobian matrix of \( f(y) \). If \( A \) is Henselian and if \( f(0) = 0 \) and \( J \notin m_A \), \( \frac{A[y]}{(y)} \), then there exists \( \tilde{y} \in m_A^r \) such that \( f(\tilde{y}) = 0 \).

Example 3.8.14. — The ring of germs of \( C^\infty \) function at the origin of \( \mathbb{R}^n \) is a Henselian local ring but it is not Noetherian. The ring of germ of analytic functions at the origin of \( \mathbb{C}^n \) is a Noetherian Henselian local ring; it is the ring of convergent power series.

Example 3.8.15. — If \( A = \mathbb{k}[\!\! [x_1, ..., x_n] \!\!] \) for some Weierstrass system over \( \mathbb{k} \), then \( A \) is a Henselian local ring by Proposition 3.8.12. Indeed, let \( P(y) \in A[y] \) satisfies \( P(0) = 0 \) and \( \frac{\partial P}{\partial y}(0) \notin (p, x) \). Thus \( P(y) \) contains a nonzero term of the form \( cy \), \( c \in \mathbb{k}^* \). Then we have \( y = P(y)Q(y) + R \) where \( R \in m_A \). Clearly \( Q(y) \) is a unit, thus \( P(R) = 0 \).

Proposition 3.8.16 (Hensel Lemma). — Let \( (A, m_A) \) be a local ring. Then \( A \) is Henselian if and only if for any monic polynomial \( P(y) \in A[y] \) such that \( P(y) = f(y)g(y) \mod m_A \) for some \( f(y), g(y) \in A[y] \) which are coprime modulo
m_A, there exists \( \tilde{f}(y), \tilde{g}(y) \in A[y] \) such that \( P(y) = \tilde{f}(y)\tilde{g}(y) \) and \( \tilde{f}(y) - f(y), \tilde{g}(y) - g(y) \in m_A[y] \).

**Proof.** — Let us prove the sufficiency of the condition. Let \( P(y) \in A[y] \) and \( a \in A \) satisfy \( P(a) \in m_A \) and \( \frac{\partial P}{\partial y}(a) \notin m_A \). This means that \( P(X) = (X - a)Q(X) \) where \( X - a \) and \( Q(X) \) are coprime modulo \( m \). Then this factorization lifts to \( A[X] \), this means \( \tilde{y} \in m_A \) such that \( P(a + \tilde{y}) = 0 \). This proves that \( A \) is Henselian.

To prove that the condition is necessary, let \( P(y) \in A[y] \) be a monic polynomial, \( P(y) = y^d + a_1 y^{d-1} + \cdots + a_d \). Let \( k := \frac{A}{m_A} \) be the residue field of \( A \), an any \( a \in A \), let us write \( \overline{a} \) for the image of \( a \) in \( k \). Let us assume that \( \overline{P}(y) = f(y)g(y) \mod m_A \) for some \( f(y), g(y) \in k[y] \) which are coprime in \( k[y] \). Let us write

\[
\begin{align*}
f(y) &= y^{d_1} + b_1 y^{d_1-1} + \cdots + b_{d_1}, \\
g(y) &= y^{d_2} + c_1 y^{d_2-1} + \cdots + c_{d_2}
\end{align*}
\]

where \( b = (b_1, \ldots, b_{d_1}) \in k^{d_1} \), \( c = (c_1, \ldots, c_{d_2}) \in k^{d_2} \). The product of polynomials \( \overline{P} = fg \) defines a map \( \Phi: k^{d_1} \times k^{d_2} \to k^d \), that is polynomial in \( b \) and \( c \) with integer coefficients, and \( \Phi(b, c) = \overline{a} := (\overline{a}_1, \ldots, \overline{a}_d) \). The determinant of the Jacobian matrix \( \frac{\partial \Phi}{\partial (b, c)} \) is the resultant of \( f(y) \) and \( g(y) \), and hence is nonzero at \((b, c)\). Using the Implicit Function Theorem [3.8.13] there exist \( \tilde{b} \in A^{d_1} \), \( \tilde{c} \in A^{d_2} \) such that \( P(y) = P_1(y)P_2(y) \) where \( P_1(y) = y^{d_1} + \tilde{b}_1 y^{d_1-1} + \cdots + \tilde{b}_{d_1} \) and \( P_2(y) = y^{d_2} + \tilde{c}_1 y^{d_2-1} + \cdots + \tilde{c}_{d_2} \).

**Proposition 3.8.17.** — ([GrDi67] 18-7-6) If \( A \) is an excellent local ring, then its Henselization \( A^h \) is also an excellent local ring.
CHAPTER 4

PRESENTATION OF OUR WORK

4.1. Study of the Artin function of polynomial equations with coefficients in the ring of power series in two variables

The aim of this part is to present our results concerning the problem of obtaining effective bounds on Artin functions of polynomials with coefficients in the ring of polynomials in one or two polynomials. The work follows on from our work done during our PhD thesis.

As seen in the previous part, M. Artin has shown the following result (cf. Theorem 3.3.12):

**Theorem 4.1.1.** — [Ar69] Let $k$ be a field and $f_1, \ldots, f_p \in k[t_1, \ldots, t_m, X_1, \ldots, X_n]$ be polynomials. Then for any integer $c \in \mathbb{N}$ there exists an integer $\beta(c) \in \mathbb{N}$ satisfying the following property: for any $\overline{x} \in k[[t_1, \ldots, t_m]]^n$ such that

$$f_1(t, \overline{x}), \ldots, f_p(t, \overline{x}) \in (t)^{\beta(c)},$$

there exists $\tilde{x} \in k[[t]]^n$ such that

$$f_1(t, \tilde{x}), \ldots, f_p(t, \tilde{x}) = 0$$

and $\tilde{x} - \overline{x} \in (t)^c$.

**Definition 4.1.2.** — For any integer $c$ we denote by $\beta(c)$ the least integer satisfying the property of Theorem 4.1.1. The function $c \mapsto \beta(c)$ is called the Artin function of $f_1, \ldots, f_N$ and, in fact, it depends only on the ideal $I$ of $k[t, X]$ generated by $f_1, \ldots, f_N$. This means that if $g_1, \ldots, g_r$ generate the same ideal $I$, then the Artin function of $g_1, \ldots, g_r$ is the same as the Artin function of $f_1, \ldots, f_p$ (See Remark 3.3.20).

When $m = 1$, we call this function the Greenberg function of $f_1, \ldots, f_p$.

Indeed a few years before, M. Greenberg had given a proof of this theorem in the case $m = 1$ (cf. Theorem 3.3.1 in the preceding part) and in this case
he showed that we may choose $\beta$ to be affine (cf. [Gre66]). This means that any Greenberg function is bounded by an affine function.

In fact M. Artin has shown a stronger result than the one stated here: $\beta(c)$ may be chosen in such a way that it depends only on the degrees of $f_1, \ldots, f_p$ and on the number of variables $n$ and $m$ (see Theorem 3.3.12 for a precise statement).

In [Ar70], M. Artin raised the problem of having estimates for the Artin function of a given system of equations. We know (see Example 3.3.3 or [Ron05b] and [Ron06a]) that for $m \geq 2$ the Artin function of a system of equations is not bounded by an affine function in general, unlike the case $m = 1$ proven by M. Greenberg. But while we know few precise results in the case $m = 1$ (as shown in Part 3.3.1), there is no general known bound for $m \geq 2$, the only known result is that such an Artin function is bounded by a computable function (see [Las78] or [BDLvdD79]). Moreover the only examples of Artin functions whose a bound is known are always very particular examples and the bounds are always affine. (see Part 3.3.4, [Ron06a] or [Di07]). Thus the question asked by M. Artin is still widely open.

The difficulty to obtain an "effective" bound in Theorem 4.1.1 can be explained by sketching the proof of this result. There exist several analogues of Theorem 4.1.1 in different situations, but there exist essentially two proofs. The first one is due to M. Artin and is based on an induction on the number $m$ of variables $t$. The second one uses model theoretical methods (ultraproducts) as presented in Part 3.3.3 and is absolutely not effective.

The proof due to M. Artin consists in several steps: first, if $I$ is the ideal of $k[t, X]$ generated by $f_1, \ldots, f_p$ and if $I = Q_1 \cap \cdots \cap Q_s$ is a primary decomposition of $I$ then the Artin function of $I$ is bounded by the sum of the Artin functions of the $Q_i$'s. Then if $Q$ is a primary ideal and $e$ is an integer such that $\sqrt{Q}^e \subset Q$, the Artin function of $Q$ is bounded by $e$ times the Artin function of $\sqrt{Q}$. Therefore we may assume that $I$ is a prime ideal. Then, if $I$ is prime, either we can apply the Implicit function theorem and the Weierstrass division Theorem to reduce the problem to the case of $m - 1$ variables $t$, either we replace $I$ by $I + (\delta)$ where $\delta$ is a well chosen minor of the Jacobian matrix of $I$ and we increase the height of the ideal (according to the Jacobian criterion since $I$ is prime). Thus we do a double induction on the height of $I$ and the number $m$ of variables $t$. But, at each step of the induction on the height of $I$ we need to replace $I$ by one of its minimal prime ideals. This double induction makes the effectivity of the Artin functions involved quite difficult to control. The only thing that can be done is to remark (as M. Artin did) that we have a control on the degrees of the generators of each ideal involved in this double induction. Hence the idea is to try to prove the existence of $\beta$ satisfying
Theorem [4.1.1] but depending only on the degrees of the generators of $I$. The control we have on the degrees of the generators of each ideal appearing in the induction is quite tedious since we need to replace at each step the ideal $I$ by one of its minimal prime. Nevertheless, for $m$ small enough (in order to avoid a heavy induction on the number $m$ of variables $t$), more precisely for $m \leq 2$, this idea gives us few effective results. This is the aim of this first part.

The first result we obtain is the following (in the case $m = 1$):

**Theorem 4.1.3.** — [Ron10a] Let $k$ be a perfect field. For any $n$, $d \in \mathbb{N}$, there exists $\beta : \mathbb{N} \rightarrow \mathbb{N}$ such that for any ideal $I$ in $k[t, X]$, with $X = (X_1, \ldots, X_n)$, generated by polynomials of degree less than $d$, and for any $c \in \mathbb{N}$ and any $x(t) \in k[[t]]^n$ such that $f(t, x(t)) \in (t)^{\beta(c)}$ for all $f \in I$, there exists $\pi(t) \in k[[t]]^n$ such that $f(t, \pi(t)) = 0$ for all $f \in I$ and $x(t) - \pi(t) \in (t)^c$. Moreover $\beta$ may be chosen to be affine, of the form $c \mapsto a(n,d)(c + 1)$ where $a(n,d)$ is bounded by a polynomial function in $d$ of degree exponential in $n$.

The new thing here, compared to Greenberg’s result, is that the affine bound of the Greenberg function is uniform in $d$ and $m$, but moreover we have a bound on the coefficients of this affine function. The fact that the bound is doubly exponential in $n$ comes from the fact that the bound on the degrees of the generators of an associated prime of $I$ is doubly exponential in $n$ (cf. [Se74] and [Te90]).

Let us remark that this result is stated for a characteristic zero field in [Ron10a]. In fact the only difficulty appearing in positive characteristic is the use of the Jacobian criterion, but this difficulty is avoided by assuming that the field is perfect.

Then we can use this result to find bounds of the Artin function of polynomials equations with coefficients in $k[t_1, t_2]$. The first result we obtain concerns the case of binomial equations. The will explain this through the following example:

Let us consider the polynomial $X^2 - Y^3$ seen as a polynomial of $k[[t_1, t_2]][X, Y]$. Let $x(t)$ and $y(t) \in k[[t_1, t_2]]$ be two given non zero formal power series. Let us denote respectively by $r$ and $s$ their vanishing order at 0. After a linear change of coordinates in $t_1$ and $t_2$ we may assume that $x(t)$ and $y(t)$ are $t_2$-regular of order $r$ and $s$, i.e. $x(0, t_2) = at_2^r$ et $y(0, t_2) = bt_2^s$ with $a$ and $b$ two non zero elements of $k$. By the Weierstrass preparation Theorem for formal power series we can write

$$x(t) = u(t) \left( t_2^r + a_1(t_1)t_2^{r-1} + \cdots + a_r(t_1) \right),$$
$y(t) = v(t) \left( t_2^{s} + b_1(t_1)t_2^{s-1} + \cdots + b_s(t_1) \right)$

where $u(t)$ and $v(t)$ are units of $\mathbb{k}[t]$ and the $a_i(t_1)$'s and $b_j(t_1)$'s are formal power series in one variable. Then let us denote by $P(t_2)$ the polynomial $x(t)u(t)^{-1}$ and by $Q(t_2)$ the polynomial $y(t)v(t)^{-1}$. These are polynomials of $\mathbb{k}[t_1][t_2]$.

Let us assume that

$$x(t)^2 - y(t)^3 \in (t)^c$$

where $c$ is a sufficiently large integer (let us say larger than $r$ and $s$). This can be rewritten as the following:

$$u(t)^2 P(t_2)^2 - v(t)^3 Q(t_2)^3 \in (t)^c.$$

Let us write

$$P(t_2)^2 = t_2^{2r} + c_1(t_1)t_2^{2r-1} + \cdots + c_{2r}(t_1)$$

$$Q(t_2)^3 = t_2^{3s} + d_1(t_1)t_2^{3s-1} + \cdots + d_{3s}(t_1).$$

Thus we can express the $c_i$'s and $d_j$'s as polynomials in the coefficients of $P(t_2)$ and $Q(t_2)$:

$$c_i = C_i(a_1, \ldots, a_r) \quad \text{and} \quad d_j = D_j(b_1, \ldots, b_s).$$

For example $C_{2r} = a_r^2$ and $C_{2r-1} = 2a_r a_{r-1}$. Then we can show (cf. lemme 4.1 Ron10a) that necessarily $2r = 3s$ and

$$u(t)^2 - v(t)^3 \in (t)^{c-2r},$$

$$c_i(t_1) - d_i(t_1) \in (t_1)^{c-2r+i} \quad \text{for any} \ i.$$

Thus we have obtained a new system of equations, formed by the equations (10) and (11) that we can consider separately since they do depend on separated sets of unknowns.

Equation (10) may be solved easily since $u(0)$ and $v(0)$ are different from zero and the point $(u(0), v(0)) \in \mathbb{k}^2$ is not in the singular locus of $X^2 - Y^3 = 0$.

Now, it is easy to check, using the Implicit Function Theorem, that the Artin function of smooth systems of equations is the identity function (see Remark 3.3.21). Hence there exist two power series $\pi(t)$ and $\tau(t) \in \mathbb{k}[t]$ such that

$$\pi(t)^2 - \tau(t)^3 = 0$$

$$\pi(t) - u(t), \tau(t) - v(t) \in (t)^c.$$

The system of equations (11) may be written as follows:

$$C_i(a_1(t_1), \ldots, a_r(t_1)) - D_i(b_1(t_1), \ldots, b_s(t_1)) \in (t_1)^c \quad \forall i.$$

This is a system of equations whose approximated solutions are power series in one variable $t_1$. Thus we may apply Theorem 4.1.3 since we know that the
degree of the polynomials $C_i - D_i$ is less than 3. Thus we know that there exist $\pi_j(t_1)$ and $\overline{b_k}(t_1)$ such that

\[ C_i(\pi_1(t_1), \ldots, \pi_r(t_1)) - D_i(\overline{b_1}(t_1), \ldots, \overline{b_s}(t_1)) = 0 \quad \forall i \]

and $\pi_j(t_1) - a_j(t_1)$, $\overline{b_k}(t_1) - b_k(t_1) \in (t_1)^{c'}$ where $c'$ can be determined in function of $c$. Hence we set

\[ \pi(t) := \pi(t) \left( t_2^2 + \pi_1(t_1) t_2^{r-1} + \cdots + \pi_r(t_1) \right), \]

\[ \overline{y}(t) = \overline{y}(t) \left( t_2^2 + \overline{b_1}(t_1) t_2^{2-1} + \cdots + \overline{b_s}(t_1) \right). \]

We will have $\pi(t)^2 - \overline{y}(t)^2 = 0$ and $\pi(t) - x(t)$, $\overline{y}(t) - y(t) \in (t)^{c'}$. We do not give here more details but this method can be used for any system of binomial equations and yields the following result:

**Theorem 4.1.4.** Let $k$ be a perfect field. Let us set $t = (t_1, t_2)$. Then the followings are satisfied:

i) For any $d, d' \in \mathbb{N}$, there exists $a_{d,d'} > 0$ satisfying the following property:

Let $I$ be a binomial ideal of $k[X_1, \ldots, X_n]$ generated by binomials $f_1, \ldots, f_p$ of degree less than $d'$. Let $c \in \mathbb{N}$ and $x_1(t), \ldots, x_n(t) \in k[t]$ satisfy $\text{ord}(x_j(t)) \leq d$ and $f_k(x_j(t)) \in (t)^{a_{d,d'}(c+1)}$ for any $j$ and $k$. Then there exists $\overline{x}_j(t) \in k[t]$ such that $f_k(\overline{x}_j(t)) = 0$ for all $k$ and $\overline{x}_j(t) - x_j(t) \in (t)^{c}$ for all $j$.

ii) For any $d' \in \mathbb{N}$ there exists a doubly exponential function in $c$, denoted by $\beta_{d'}$, such that for any binomial ideal $I$ of $k[X_1, \ldots, X_n]$ generated by binomials of degree less than $d'$, the Artin function of $I[k[t]]/X$ is bounded by $\beta_{d'}$.

Once more in [Ron10a] this result is stated for algebraically closed fields of characteristic zero, but we can extend to algebraically closed fields of positive characteristic since it is proven using Theorem 4.1.3. In [Ron10a] the fact that $k$ is algebraically closed is used to reduce the problem to the case of a prime ideal. Indeed if $I$ is a binomial ideal and $k$ is algebraically closed then the minimal primes of $I$ are binomial ideals. Nevertheless it is enough to replace $I$ by its radical that is always a binomial ideal (see [Esi96]). Indeed, by slightly modifying the previous example, if $I = ((X^2 - Y^3)^2)$ then $(u(0), v(0))$ is in the singular locus of $(X^2 - Y^3)^2 = 0$ which is the entire cusp and we cannot apply Remark 3.3.21. Thus we need to replace $I$ by its radical. If $I$ is radical then the singular locus of $V(I)$ is included in the union of the coordinates hyperplanes, thus $(u(0), v(0))$ is never included in the singular locus of $V(I)$.

Let us remark here that the only two known examples of Artin functions which are not bounded by affine functions are Artin functions of binomial equations.
In these two cases we just know that these Artin functions are bounded from below by a polynomial function of degree 2 (cf. Example 3.3.33 or \cite{Ron05b} and \cite{Ron06a}).

More generally we can use Theorem \[4.1.3\] to study Artin functions of systems of polynomial equations with coefficients in \(k\left[t_1,t_2\right]\). The idea is to follow the proof of M. Artin of Theorem \[4.1.1\] in the case \(m = 2\). Hence we obtain the following result:

**Theorem 4.1.5.** — \[Ron13a\] Let \(k\) be a perfect field and let \(f_1,\ldots, f_p\) be polynomials of \(k[t_1,t_2][X_1,\ldots,X_n]\) generating an ideal \(I\). Let us denote by \(H\) an ideal of \(k[t,X]\) defining the non-smoothness locus of the morphism \(k[t_1,t_2] \rightarrow k[t_1,t_2,X_1,\ldots,X_n]_{(f_1,\ldots,f_p)}\), i.e. the critical locus of the projection \(\pi: V(f_1,\ldots,f_p) \subset k^2 \times k^n \rightarrow k^2\). Then for any \(d \in \mathbb{N}\) there exists \(a(d) > 0\) and \(b(d) > 0\) such that for any \(x(t) \in k[t]^n\) satisfying

\[
  f(t,x(t)) \in (t)^{a(d)(c+1)} \quad \text{for all} \quad f \in I
\]

and \(\exists h \in H, \quad h(t,x(t)) \notin (t)^{b(d)}\), there exists \(\pi(t) \in k[[t]]^n\) such that

\[
  f(t,x(t)) = 0 \quad \text{for all} \quad f \in I
\]

and \(\pi(t) - x(t) \in (t)^c\).

Let us remark that if \(I\) is generated by polynomials of \(k[X_1,\ldots,X_n]\) (i.e. the \(f_i\)'s do not depend on \(t_1\) and \(t_2\)), then \(H\) is just an ideal defining the singular locus of \(f = 0\) in \(k^n\).

**Remark 4.1.6.** — There are several ways to define such an ideal \(H\). The first definition is due to Elkik \[Elk73\] and it is the following:

Let \(A\) be a Noetherian ring (here \(A = k[[t_1,t_2]]\)) and let \(f_1,\ldots, f_p \in A[X_1,\ldots,X_n]\). Let \(E\) be a subset of \(\{1,\ldots,p\}\) whose cardinal is \(h\) for some integer \(h\). We denote by \(\Delta_E(f)\) the ideal of \(A[X_1,\ldots,X_n]\) generated by the \(h \times h\) minors of the Jacobian matrix \(\left(\frac{\partial f_i}{\partial X_j}\right)_{i \in E, 1 \leq j \leq n}\). (This ideal is zero if \(h > n\)). We define the following ideal of \(A[X_1,\ldots,X_n]\):

\[
  H_{f_1,\ldots,f_p} := \sum_E \Delta_E(f)((f_i, i \in E) : I)
\]

where the sum runs over all subsets \(E\) of \(\{1,\ldots,p\}\). This ideal \(H_{f_1,\ldots,f_p}\) defines the non-smoothness locus of the morphism \(A \rightarrow A[X_1,\ldots,X_n]_{(f_1,\ldots,f_p)}\). I have never found any reference about the fact that this ideal is independent of the choice of the generators \(f_1,\ldots, f_p\) of the ideal \(I\). Most of the time nothing is said about
4.1. ARTIN FUNCTION

This problem and sometimes it is claimed that it is easy to check that it does not depend on this choice without giving a proof of it. Unfortunately, as shown on an example in [Ron13a] (we consider $I = (X_1, X_2) \cap (X_3, X_4)$ defining a non complete intersection singularity), this definition depends on the choice of the generators of $I$. On the other hand the radical of this ideal $H_{f_1, \ldots, f_p}$ does not depend on the generators of $I$. In fact it does not depend on the choice of a representation of the morphism $A \to A[X_1, \ldots, X_n]^{\langle f_1, \ldots, f_p \rangle}$. In the previous theorem, we can take any ideal $H$ whose radical is equal to the radical of $H_{f_1, \ldots, f_p}$.

This theorem may be rephrased as saying that the Artin function of a system of equations in $k[t_1, t_2][X_1, \ldots, X_n]$ is bounded by an affine function if we consider approximated solutions whose contact order with the critical locus of $\pi$ is bounded.

In fact the result proven in [Ron13a] is more precise and may be stated more easily by using the norm induced by ord on $k J_{t_1, t_2}$ as follows (see Remark 3.2.3 and Remark 3.3.6 in the previous part for a precise definition of this topology and for the relation with Łojasiewicz inequalities):

**Theorem 4.1.7.** — [Ron13a] Let $A := k[t_1, t_2]$ where $k$ is a perfect field. Then there exist constants $K_1, K_2 > 0$ such that for any $d \geq 2$ and $n \geq 1$, for any ideal $I = (f_1, \ldots, f_p)$ of $k[t_1, t_2, X_1, \ldots, X_n]$ generated by polynomials of degrees less than $d$ with $V(I) \neq \emptyset$, we have the following inequality:

$$||f(\pi)|| \geq (K_1 d(\pi, f^{-1}(0)))^{d \left( \frac{1}{n!} \right)^{K_2 n}} \forall \pi \in A^n \setminus V(H)$$

where $H$ is an ideal of $k[t_1, t_2, X_1, \ldots, X_m]$ defining the critical locus of the projection $V(f_1, \ldots, f_p) \subset k^2 \times k^n \to k^2$.

Let us mention that this result is no more valid when $m \geq 3$ (i.e. $A = k[t_1, \ldots, t_m]$ with $m \geq 3$). Indeed let us consider the following example (see [Ron06a]):

**Example 4.1.8.** — Let $A := k[t_1, t_2, t_3]$ and let $f := X_1 X_2 - X_3 X_4$. Here we can take for $H$ the ideal $(X_1, X_2, X_3, X_4)$. For $c \geq 3$ let us set

$$\pi_1^{(c)} := t_1^c, \quad \pi_2^{(c)} := t_2^c, \quad \pi_3^{(c)} := t_1 t_2 - t_3^c.$$

Then there exists $\pi_4^{(c)} \in A$ such that $\pi_1^{(c)} \pi_2^{(c)} - \pi_3^{(c)} \pi_4^{(c)} \in (t)^{c^2}$. This means that $||f(\pi^{(c)})|| \leq e^{-c^2}$ for any $c \geq 3$. Moreover it is shown in [Ron06a] that
any solution $\tilde{x} \in A^4$ of $f = 0$ satisfies
\[
\min_{i=1,...,4} \{\text{ord}(\pi_i^{(c)} - \tilde{x}_i)\} \leq c.
\]
this means that $d(\pi_i^{(c)}, f^{-1}(0)) \geq e^{-c}$ for all $c \geq 3$. Thus we see that there do
not exist constants $a > 0$ and $b > 0$ such that $||f(\pi^{(c)})|| \geq a \cdot d(\pi^{(c)}, f^{-1}(0))^b$
for all $c \geq 3$, but $||H(\pi^{(c)})|| = e^{-2}$ is a constant not depending on $c$.

From Theorem 4.1.5 we can deduce the following corollary which is an ana-
logue of the second part of Theorem 4.1.4 for isolated singularities:

**Corollary 4.1.9.** — [Ron13a] We use the notation of the previous theorem. Let $I$ be an ideal of $k[t_1, t_2, X_1, ..., X_n]$ such that the critical locus of the projec-
tion $V(f_1, ..., f_p) \subset k^2 \times k^n \rightarrow k^2$ has an isolated singularity. Then the Artin
function of $I$ is bounded by a doubly exponential function.

We can mention here that Hickel has given a bound of the Artin function
of ideals $I$ of $k[t, X_1, ..., X_n]$ ($t$ being a single variable) satisfying the same
hypothesis (see Remark 3.3.10 ii) or [Hic93]).

### 4.2. Morphisms of local algebras and nested approximation

In this part we present our results concerning some regularity properties of
morphisms between rings of formal or convergent power series. This is related
to a particular case of Artin approximation with constraints. We begin with
a short survey of the problem. This work has been motivated by a question
that S. Izumi asked us in 2005 during a stay at Osaka. The problem raised by
S. Izumi was the possibility to extend in positive characteristic some known
results in characteristic zero. The question was more specifically related to
the Chevalley function of a morphism between local rings. One of our main
contributions has consised to extend known results in characteristic zero to the
positive characteristic case (see Theorem 4.2.16 in particular).

#### 4.2.1. Persistence of properties of an analytic morphism after com-
pleation.

Let us begin by stating two classic corollaries of the Weierstrass
preparation Theorem (see [Ho61], [Mal68] or [To72] for example):

**Proposition 4.2.1.** — Let $k$ be a valued field. Let $\varphi : A \rightarrow B$ be a morphism
of analytic algebras. Then $B$ is finite over $A$ if and only if $\hat{B}$ is finite over $\hat{A}$.

**Proof.** — Indeed by the Weierstrass preparation Theorem, $B$ is finite over $A$
if and only if $B/\mathfrak{m}_B \simeq B/\mathfrak{m}_B$ and $\hat{A}/\mathfrak{m}_A \simeq A/\mathfrak{m}_A$. By applying the Weierstrass preparation Theorem of algebras of formal
power series we obtain the result. □
Proposition 4.2.2. — Let \( \varphi : A \to B \) be a morphism of analytic algebras. Then

i) \( \varphi \) is surjective if and only if \( \hat{\varphi} \) is surjective.

ii) If \( \hat{\varphi} \) is injective then \( \varphi \) is injective.

iii) \( \varphi \) is an isomorphism if and only if \( \hat{\varphi} \) is an isomorphism.

Proof. — First of all if \( \varphi \) or \( \hat{\varphi} \) is surjective then \( B \) is finite over \( A \) or \( \hat{B} \) is finite over \( \hat{A} \), thus \( B \) is finite over \( A \) and \( \hat{B} \) is finite over \( \hat{A} \) by the previous proposition. Since \( B/\mathfrak{m}_B \cong \hat{B}/\mathfrak{m}_{\hat{B}} \), by using Nakayama’s Lemma we obtain the first statement.

If \( \hat{\varphi} \) is injective then \( \varphi \) is obviously injective.

If \( \varphi \) is an isomorphism then \( \hat{\varphi}^{-1} \) is the inverse of \( \hat{\varphi} \) which is also an isomorphism. If \( \hat{\varphi} \) is an isomorphism then \( \varphi \) is injective and surjective by the preceding points, hence \( \varphi \) is an isomorphism.

In [Gro60], A. Grothendieck conjectured that an injective morphism of analytic algebras induces an injective morphism between the completions. More generally we can ask if an injective morphism of analytic algebras \( \varphi : A \to B \) satisfies \( \hat{\varphi}^{-1}(B) = A \). A similar question is the following: if the image of a formal power series is convergent is it the image of a convergent power series, i.e. if \( i \varphi : A \to B \) is a morphism of analytic algebras, do we have \( \hat{\varphi}(A) \cap B = \varphi(A) \)? The first answer to these questions appeared in the paper [Gab71] by A. Gabrielov and is negative (see Examples 3.5.12 and 3.5.14 of the previous part for a detailed presentation of this example).

4.2.2. Chevalley function. — In the paper [Ch43] the following result is proven:

Theorem 4.2.3 (Chevalley’s Lemma). — Let \( A \) be a complete local ring with maximal ideal \( \mathfrak{m}_A \). Let \( (I_n) \) be a decreasing sequence of ideals of \( A \) such that \( \cap_n I_n = \{0\} \). Then there exists a function \( \beta : \mathbb{N} \to \mathbb{N} \) such that \( I_{\beta(n)} \subset \mathfrak{m}_A^n \) for any positive integer \( n \).

Thus we deduce the following result (by applying the previous theorem to the sequence \( (\varphi^{-1}(\mathfrak{m}_B^n) )_n \) of ideals of \( A/\text{Ker}(\hat{\varphi}) \)):

Corollary 4.2.4. — Let \( \varphi : A \to B \) be a morphism of local complete \( k \)-algebras. Then there exists a function \( \lambda : \mathbb{N} \to \mathbb{N} \) such that

\[
\forall n \in \mathbb{N}, \quad \varphi^{-1}(\mathfrak{m}_B^{\lambda(n)}) \subset \text{Ker}(\varphi) + \mathfrak{m}_A^n.
\]

The least function \( \lambda \) satisfying this property is called the Chevalley function of the morphism \( \varphi \). This is an increasing function.
A natural question is to know the behaviour or the growth of this function. In particular what are the morphisms whose Chevalley function is bounded by an affine function?

**Example 4.2.5.** — [Ron09b] Once more, as the example of Gabrielov, this example is inspired by Osgood’s Example (see Example 3.5.12) and a remark of Abhyankar [Ab56].

Let $\alpha : \mathbb{N} \to \mathbb{N}$ be an increasing function and let $\mathbb{k}$ be a valued field. Let $(n_i)_i$ be a sequence of natural numbers such that $n_{i+1} > \alpha(n_i)$ for any $i$ and such that the element $\xi(Y) := \sum_{i \geq 1} Y^{n_i}$ is transcendental over $\mathbb{k}(Y)$ (such an element exists according to the constructive proof of Lemma 1 in [McLSc39]).

Let us define the morphism $\varphi : \mathbb{k}\{x_1, x_2, x_3\} \to \mathbb{k}\{y_1, y_2\}$ by

$$((\varphi(x_1), \varphi(x_2), \varphi(x_3)) = (y_1, y_1y_2, y_1\xi(y_2)).$$

We show exactly as for Osgood’s Example (see Example 3.5.12) that $\hat{\varphi}$ is injective. For any positive natural number $i$ we define:

$$f_i := x_1^{n_i-1} x_3 - \left( x_2^{n_1} x_1^{n_{i-1}} - \cdots + x_2^{n_i-1} x_1^{n_{i-1}} - x_2^{n_i} \right),$$

Then we get:

$$\varphi(f_i) = y_1^{n_i} \xi(y_2) - y_1^{n_i} \sum_{k=1}^{i} y_2^{n_{k}} \in m_{A}^{n_i+1} \subset m_{B}^{\alpha(n_i)}$$

But $f_i \notin m_{A}^{n_i+1}$ thus $\beta(n_i + 1) > \alpha(n_i)$ where $\beta$ is the Chevalley function associated to $\varphi$. Because $n_i \to +\infty$ when $i \to +\infty$, we get $\limsup \frac{\beta(n_i)}{\alpha(n_i)} \geq 1$.

Thus for any increasing function $\alpha$ there exists a morphism of analytic algebras whose Chevalley function increases faster than $\alpha$. In particular there exist such morphisms whose Chevalley function is not computable since there exist increasing functions whose growth is larger than any computable function.

Once more this shows that the question has no good general answer.

### 4.2.3. Relation with the nested Artin approximation property.

This part has been inspired by the work of Becker (see Part 0 of [Be77a]).

Let $\varphi : \mathbb{k}\{x_1, \ldots, x_n\}/I \to \mathbb{k}\{y_1, \ldots, y_m\}$ be an injective morphism of analytic algebras. Let us denote by $\varphi' : \mathbb{k}\{x_1, \ldots, x_n\} \to \mathbb{k}\{y_1, \ldots, y_m\}$ the induced morphism. Let us denote by $x$ and $y$ the multi-variables $(x_1, \ldots, x_n)$ and $(y_1, \ldots, y_m)$. Then $\hat{\varphi}$ is injective if and only $\text{Ker}(\varphi') = I \mathbb{k}[x] = \text{Ker}(\varphi') \mathbb{k}[x]$. Thus in order to study the first question we may assume that $\varphi$ is a morphism of convergent power series rings, $\varphi : \mathbb{k}\{x\} \to \mathbb{k}\{y\}$, and investigate under what conditions $\text{Ker}(\hat{\varphi})$ is generated by $\text{Ker}(\varphi)$.

We can remark that $\text{Ker}(\varphi) \mathbb{k}[x] = \text{Ker}(\varphi) := \bigcap_{c \in \mathbb{N}} (\text{Ker}(\varphi) + (x)^c \mathbb{k}[x])$ is the
closure of $\ker(\varphi)$ in $k[[x]]$ for the Krull topology of this ring.

Now let us consider the proof of Lemma 3.5.1:

We denote by $\varphi_i(y) \in k\{y\}$ the image of $x_i$ under $\varphi$ for $1 \leq i \leq n$. Let us consider $\overline{f} \in \ker(\hat{\varphi})$. The Taylor expansion of $\overline{f}(x)$ yields

$$\overline{f}(x) = \overline{f}(x) - \overline{f}(\varphi(y)) = \sum_{\alpha \in \mathbb{N}^n, \alpha \neq 0} \frac{1}{\alpha_1! \cdots \alpha_n!} \overline{f}(\varphi(y))(x_1 - \varphi_1(y))^{\alpha_1} \cdots (x_n - \varphi_n(y))^{\alpha_n}.$$ 

Thus there exist $\overline{g}_i \in k[[x,y]]$, for $1 \leq i \leq n$, such that

$$\overline{f}(x) + \sum_{i=1}^n (x_i - \varphi_i(y)) \overline{g}_i(x,y) = 0.$$ 

Then $\overline{f} \in \overline{\ker(\varphi)}$ if and only if for any $c \in \mathbb{N}$ there exists $f_c \in \ker(\varphi)$ such that $\overline{f}(x) - f_c(x) \in (x)^c$. The Taylor expansion of $f_c$ shows us that $\overline{f} \in \ker(\varphi)$ if and only if there exist $f_c(x) \in k\{x\}$ and $g_{i,c}(x,y) \in k\{x,y\}$ such that

$$f_c(x) + \sum_{i=1}^n (x_i - \varphi_i(y)) g_{i,c}(x,y) = 0 \quad \forall c \in \mathbb{N}$$

and $f_c(x) - \overline{f}(x) \in (x)^c \quad \forall c \in \mathbb{N}$.

Now let us fix $P(F_1, G_1) \in k\{x,y\}[F_1, \ldots, F_r, G_1, \ldots, G_s]$. We say that $P$ satisfies the nested Artin approximation property if:

$$\forall \overline{f} \in k[[x]]^r, \forall \overline{g} \in k[[x,y]]^s, \text{ such that } P(\overline{f}, \overline{g}) = 0, \forall c \in \mathbb{N},$$

$$\exists f \in k\{x\}^r, \exists g \in k\{x,y\}^s \text{ such that } P(f, g) = 0,$$

and $\overline{f}_i - f_i \in (x)^c, \overline{g}_j - g_j \in (x,y)^c$ for all $1 \leq i \leq r, 1 \leq j \leq s$.

Thus we deduce the following proposition form the remark at the begin of this part (in fact we have shown here something weaker; for the complete proof of this proposition see [Ron08a]):

**Proposition 4.2.6.** — [Ron08a] Let us consider the following equation:

(E1) $$F(x) + \sum_{i=1}^n (x_i - \varphi_i(y)) G_i(x,y) = 0$$

Thus $\ker(\hat{\varphi}) = \ker(\varphi)k[[x]]$ if and only if Equation (E1) satisfies the nested Artin approximation property.

We can generalize this result:

**Proposition 4.2.7.** — [Ron08a] Let us consider the following equation:

(E2) $$F(x) + \sum_{i=1}^n (x_i - \varphi_i(y)) G_i(x,y) + h(y) = 0$$
where \( h \in \mathbb{k}\{y\} \). Then Equation (E2) satisfies the nested Artin approximation property if and only if \( h \in \varphi(\mathbb{k}\{x\}) \).

In particular Equation (E2) satisfies the nested Artin approximation property for any \( h \) convergent if and only if
\[
\hat{\varphi}(\mathbb{k}[x]) \cap \mathbb{k}\{y\} = \varphi(\mathbb{k}\{x\}).
\]

We also obtain the following result which is an example of positive answer to Problem 2 in Part 3.5:

**Proposition 4.2.8.** — [Ron08a] Let us consider the following equation:

\[
(E2) \quad P(F(x),G(x,y)) := F(x) + \sum_{i=1}^{n} (x_i - \varphi_i(y))G_i(x,y) + h(y) = 0.
\]

Then there exists a function \( \beta : \mathbb{N} \rightarrow \mathbb{N} \) satisfying the following property:

\[
\forall c \in \mathbb{N}, \forall \overline{f} \in \mathbb{k}[x]^r, \forall \overline{g} \in \mathbb{k}[x,y]^s, \text{ such that } P(\overline{f}, \overline{g}) \in (x,y)^{\beta(c)},
\]

\[
\exists f \in \mathbb{k}[x]^r, \exists g \in \mathbb{k}[x,y]^s \text{ such that } P(f, g) = 0,
\]

and \( \overline{f}_i - f_i \in (x)^c, \overline{g}_j - g_j \in (x,y)^c \) for all \( 1 \leq i \leq r, 1 \leq j \leq s \).

Moreover the least function satisfying this property is exactly the Chevalley function of the morphism \( \varphi \).

**4.2.4. Definitions et main theorem.** — Nevertheless there is an important case where the previous questions have a positive or "good" answer. This is the case of regular morphisms in the sense of Gabrielov. Before defining these morphisms we have to give some preliminary definitions. The reader can consult [Ron09b] for having details about the proofs.

**Definition 4.2.9.** — Let \( \mathbb{k} \) be a valued field. Let \( \varphi : A \rightarrow B \) be a homomorphism of local \( \mathbb{k} \)-algebras and let us assume that \( \text{Gr}_{m_B}B \) is an integral domain (this is the case for example when \( B \) is regular). In this case \( \text{ord}_B \), defined by \( \text{ord}_B(f) := \max\{n \in \mathbb{N} / f \in m_B^n\} \), is a valuation. Let us still denote by \( \varphi \) the morphism induced on \( A/\text{Ker}(\varphi) \) and let us denote by \( \nu := \text{ord}_B \circ \varphi \) the valuation defined on the fraction field of \( A/\text{Ker}(\varphi) \). Let \( \mathbb{k}_\nu \) be the residue field of \( \nu \), i.e. \( \mathbb{k}_\nu := A_\nu/m_\nu \) where \( A_\nu \) is the valuation ring of \( \nu \) and \( m_\nu \) its maximal ideal, and let \( \text{tr.deg}_\mathbb{k}\nu \) denote the transcendence degree of \( \mathbb{k} \rightarrow \mathbb{k}_\nu \).

If \( \varphi(m_A) = \{0\} \), then we set \( r_1(\varphi) := 0 \); otherwise we set \( r_1(\varphi) := \text{tr.deg}_\mathbb{k}\nu + 1 \).

Moreover we define

\[
r_2(\varphi) := \dim \left( \frac{\hat{A}}{\text{Ker}(\hat{\varphi})} \right),
\]

\[
r_3(\varphi) := \dim \left( \frac{A}{\text{Ker}(\varphi)} \right).
\]
Remark 4.2.10. — We can give a characterisation of $r_1$ in the case $\text{char}(k) = 0$ which is more geometric: if $(A, m)$ is a local $k$-algebra, we denote by $\Omega^1_k(A)$ the $A$-module of the Kähler differentials and $\Omega^1_k(A) := \frac{\Omega^1_k(A)}{\cap_{i=0}^{m} \Omega^i_k(A)}$ the $A$-module of separated Kähler differentials (see [GrDi64] Part 20.7 or [Sp90a]). If $\varphi : A \rightarrow B$ is a morphism of local $k$-algebras then there exists a unique morphism of $B$-modules denoted by $\varphi^1 : \Omega^1_k(A) \otimes_A B \rightarrow \Omega^1_k(B)$ which is compatible with the canonical derivations $A \rightarrow \Omega^1_k(A)$ and $B \rightarrow \Omega^1_k(B)$. Then we have $r_1(\varphi) = \text{rank}_B(\varphi^1(\Omega^1_k(A)) \otimes_A B)$. In particular if $k = \mathbb{R}$ or $\mathbb{C}$, we denote by $\Phi : (X, 0) \rightarrow (Y, 0)$ the morphism of analytic spaces induced by $\varphi$. Then $r_1(\varphi)$ is the generic rank of the jacobian of $\Phi$, and it is equal to the dimension over $k$ of the image of $\Phi$. Moreover $r_2(\varphi)$ is the dimension of the formal Zariski closure of the image of $\Phi$ and $r_3(\varphi)$ is the dimension of the analytic Zariski closure of the image of $\Phi$.

In fact the characterization of $r_1(\varphi)$ involving the rank of the jacobian matrix is the usual definition since regular morphisms were defined only in characteristic zero before our work. Since one our goals was to extend to the positive characteristic case the study of regular morphisms we needed to extend the definition of $r_1(\varphi)$ to this setting. But the rank of the jacobian matrix of a morphism of formal power series rings may be zero even in this case although it should not be zero for our purpose. For instance the morphism $\varphi : k[[x]] \rightarrow k[[y]]$ defined by $\varphi(x_i) = y_i^p$ for any $i$ satisfies $\varphi^1((y_i)^p) \subset (x_i)^p$ for any $c \in \mathbb{N}$, i.e. its Chevalley function is bounded by a linear function. In characteristic zero this is a characterization of regular morphisms, thus this morphism has to be regular. If $\text{char}(k) = 0$, then $r_1(\varphi) = n$ but if $\text{char}(k) = p > 0$ the rank of the jacobian matrix of $\varphi$ is zero when we want $r_1(\varphi) = n$ (see Definition 4.2.13).

This why we need to define $r_1(\varphi)$ in a way that do not involve derivatives.

Remark 4.2.11. — We can also give the following interpretation of $r_1(\varphi)$ when $\varphi : k[[x]] \rightarrow k[[y]]$ is a morphism of power series rings. We can define a total order on the set of monomials in $y_1, \ldots, y_m$ as done in Example 3.1.10. Let us use the notations of Example 3.1.10. Then $r_1(\varphi)$ is the dimension of the minimal cone of $\mathbb{R}^m$ containing $\{\exp(\varphi(f)) \mid f \in k[[x]]\}$.

Lemma 4.2.12. — With the previous notation we have $r_1(\varphi) \leq r_2(\varphi) \leq r_3(\varphi)$.

Proof. — The first inequality is the Abhyankar’s inequality for the valuation $\nu$. The second one comes from the fact that $\text{ht}(\text{Ker}(\varphi)) = \text{ht}(\text{Ker}(\varphi) \tilde{A}) \leq \text{ht}(\text{Ker}(\tilde{\varphi}))$. □
Definition 4.2.13. — Let \( \varphi : A \rightarrow B \) be a morphism of analytic algebras. We say that \( \varphi \) is regular in the sense of Gabrielov if \( r_1(\varphi) = r_3(\varphi) \).

Next theorem asserts that regular morphisms are the ones for which the questions of the introduction have a positive answer:

**Theorem 4.2.14.** — Let \( k \) be a valued field and let \( \varphi : A \rightarrow B \) be a morphism of analytic algebras where \( B \) is regular. Let us consider the following properties:

1. \( r_1(\varphi) = r_2(\varphi) = r_3(\varphi) \).
2. There exist \( a \geq 1, b \geq 0 \) such that \( \hat{\varphi}^{-1}(m_B^n + b) \subset Ker(\hat{\varphi}) + m_A^n \) for all \( n \in \mathbb{N} \).
3. \( \hat{\varphi}(\hat{A}) \cap B = \varphi(A) \).

Then the following implications are satisfied:

1. \( (i) \iff (ii) \) for any valued field \( k \).
2. \( (i) \implies (iii) \) if \( k = \mathbb{R} \) or \( \mathbb{C} \), or if \( k \) is any valued field and \( A \) is regular.
3. \( (iii) \implies (i) \) if \( k = \mathbb{R} \) or \( \mathbb{C} \).

The equivalence \( (i) \iff (ii) \) remains true if \( A = k[x_1, \ldots, x_n] \) and \( B = k[y] \) where the family \( (k[x_1, \ldots, x_n])_{n \in \mathbb{N}} \) is a Weierstrass system over a field \( k \) (see Definition 3.2.17 for the definition of a Weierstrass system). Moreover the implication \( (i) \implies (iii) \) is satisfied if \( A = k[x] \) and \( B = k[y] \).

**Remark 4.2.15.** — If \( \text{char}(k) = 0 \) and if \( B \) is an integral domain, the existence of a resolution of singularities for \( B \) gives an injective morphism \( \pi : B \rightarrow k[y] \) such that \( r_1(\pi) = \dim(B) \). Thus the previous theorem remains valid if \( \text{char}(k) = 0 \) and \( B \) is an integral domain.

Historically, A. Gabrielov first proved that \( r_1(\varphi) = r_2(\varphi) \) implies \( r_1(\varphi) = r_3(\varphi) \) when \( k = \mathbb{C} \) or \( \mathbb{R} \) [Gab71] (the reverse implication is trivial) and deduced easily that \( (i) \implies (iii) \). The proof of this result of Gabrielov is quite difficult and several people tried to give a correct proof of it (see [To90] or [Sp90a] - we confess having given a wrong proof of this result). We may reformulate this implication in the following way: if \( (X, 0) \) is a germ of irreducible formal space of which a piece is the image of a germ of analytic space by an analytic map, then \( (X, 0) \) is a germ of analytic space.

The proof of \( (i) \implies (iii) \) for any valued field has been given by the author [Ron09b].

The equivalence \( (i) \iff (ii) \) has first been proven by S. Izumi for \( k = \mathbb{C} \) or \( \mathbb{R} \) [Iz86], then for any characteristic zero field \( k \) [Iz89], then by the author for any field \( k \) and for Weierstrass systems [Ron09b]. The implication \( (iii) \implies (i) \) has been proven by P. Eakin and G. Harris in the case \( A \) is regular [EaHa77].
(they asserts that they prove it for any valued field of characteristic zero, but their proof works only when $k = \mathbb{R}$ or $\mathbb{C}$), then by P. Milman [Mi78a] for any $A$.

We will only sketch the proofs of (i) $\iff$ (ii) $\implies$ (iii).

The reader may consult [BiMi82], [BiMi87], [BiMi98], [Pał92] for a study of global properties of regular morphisms and the relation with the composite functions problem in the $C^\infty$ case.

4.2.5. Monomialisation of a morphism and proof of (i) $\iff$ (ii). —

We present here a very useful result about the monomialisation of a morphism between power series rings. This one has been proven in characteristic zero in [EaHa77] then in positive characteristic in [Ron09b]. This is the key tool for the proof of (i) $\iff$ (ii) $\implies$ (iii).

**Theorem 4.2.16.** — [EaHa77] [Ron09b] Let $\varphi : k\{x\} \to k\{y\}$ be a morphism of convergent power series rings where $k$ is a valued field. Then there exist morphisms $\sigma_1 : k\{x\} \to k\{x\}$, $\sigma_2 : k\{y\} \to k\{y\}$ and $\varphi : k\{x\} \to k\{y\}$ satisfying the following properties:

i) The morphism $\sigma_1$ is the composition of $k$-automorphisms of $k\{x\}$, of morphisms $\chi_d$ ($d \in \mathbb{N}^*$ is any prime number) defined $\chi_d(x_i) = x_i^d$ and $\chi_d(x_i) = x_i$ for $i \neq 1$, and the morphism $q$ defined by $q(x_1) = x_1x_2$ and $q(x_i) = x_i$ for $i \neq 1$.

ii) The morphism $\sigma_2$ is the composition of $k$-automorphisms of $k\{y\}$ and of the morphism $q$ defined by $q(y_1) = y_1y_2$ and $q(y_i) = y_i$ for $i \neq 1$.

iii) The morphism $\varphi$ satisfies

\[ \varphi(x_i) = y_i^{\alpha_i}v_i \text{ where } v_i \text{ is a unit and } \alpha_i \in \mathbb{N}, \text{ for } i \leq r, \text{ if } \text{char}(k) = p > 0 \]

or

\[ \varphi(x_i) = y_i \text{ for } i \leq r, \text{ if } \text{char}(k) = 0 \]

and

\[ \varphi(x_{r+1}) = \cdots = \varphi(x_n) = 0. \]

Moreover $r = r_1(\varphi) = r_1(\varphi)$.

iv) The following diagram is commutative:

\[ \begin{array}{c}
\mathbb{k}\{x\} \xrightarrow{\varphi} \mathbb{k}\{y\} \\
\downarrow \sigma_1 \quad \quad \quad \downarrow \sigma_2 \\
\mathbb{k}\{x\} \xrightarrow{\varphi} \mathbb{k}\{y\}
\end{array} \]

This result remains valid if we replace $\mathbb{k}\{x\}$ and $\mathbb{k}\{y\}$ by $\mathbb{k}[x]$ and $\mathbb{k}[y]$ where the family $(\mathbb{k}[x_1, \ldots, x_i])_{i \in \mathbb{N}}$ is a Weierstrass system over a field $k$.

Now we can sketch the proof of the equivalence of (i) and (ii) in Theorem 4.2.14.
Let us prove (i) $\implies$ (ii). By replacing $A$ by $A/\ker(\varphi)$ we may assume that $\varphi$ is injective. There exists an injective and finite morphism $k\{x\} \longrightarrow A$. If we compose $\varphi$ with this morphism the ranks $r_1$, $r_2$ and $r_3$ do not change (see Lemma 2.4 [Ron09a]) and this new morphism satisfies (ii) if and only if $\varphi$ does (see Lemmas 4.4 and 4.5 [Ron09a]). Thus we may assume that $A = k\{x\}$ is regular and $\varphi$ is injective.

The specific forms of $\sigma_1$ and $\sigma_2$ allow us to show easily that if there exist constants $a$ and $b$ such that

\begin{equation}
  a \ord(f) + b \geq \ord(\overline{\varphi}(f))
\end{equation}

for all $f \in k\{x\}$, then there exist constants $a'$ and $b'$ such that

\begin{equation}
  a' \ord(f) + b' \geq \ord(\varphi(f))
\end{equation}

for all $f \in k\{x\}$ (see Lemma 3.3 [Ron09a]). If $r_1(\varphi) = n$ then $r_1(\overline{\varphi}) = n$ and the particular form of $\overline{\varphi}$ shows that $\overline{\varphi}$ satisfies an inequality of the form (13) and thus $\varphi$ satisfies (ii).

The reverse implication is shown as did Izumi in [Iz89] with the help of Theorem 4.2.16 and the use of Hilbert-Samuel functions. We do not give more details here.

**Remark 4.2.17.** — We can give an alternative proof of (i) $\implies$ (ii) in the more general case where $B$ is an integral domain. In this case we can define $r_1(\varphi)$ as follows. There exist a finite number of divisorial valuations $\mu_1, \ldots, \mu_p$ such that

\[ \ord_B(g) = \min_{i=1}^p \mu_i(g) \quad \forall g \in B. \]

These valuations are the Rees valuations of $m_B$ (see [HuSw06]). Let $\nu_i$ be the valuation defined on $A/\ker(\varphi)$ by the formula $\nu_i = \mu_i \circ \varphi$. We define $r_1(\varphi)$ as follows:

\[ r_1(\varphi) := \min_{i=1}^p \trdeg_k \nu_i + 1. \]

Then $r_2(\varphi)$ and $r_3(\varphi)$ are defined as before. Once more Abhyankar’s inequality asserts that $r_1(\varphi) \leq r_2(\varphi)$. If $r_1(\varphi) = r_2(\varphi) = \dim \left( \frac{A}{\ker(\varphi)} \right)$, then $\nu_1, \ldots, \nu_p$ are Abhyankar valuations, thus there exists a constant $a > 0$ such that $\nu_i(f) \leq a \ord_{A/\ker(\varphi)}(f)$ for $1 \leq i \leq p$ ([Sp90b] or Proposition 6.5 [Te13]). Hence $\ord_B(\varphi(f)) \leq a \ord_{A/\ker(\varphi)}(f)$ for any $f \in A$.

Nevertheless Theorem 4.2.16 is interesting since it is used to prove (i) $\implies$ (iii). Moreover it shows also that the fact that some morphisms do not satisfy the equivalent properties of Theorem 4.2.14 comes from the use of quadratic transforms in the monomialization process.
\textbf{Remark 4.2.18.} — In Theorem 4.2.16 if we can construct \( \varphi \) without using the morphism \( q : k[x] \to k[y] \), then we can see that \( \varphi \) satisfies (ii) and (iii) of Theorem 4.2.14. The fact that some morphism do not satisfy the equivalent properties of Theorem 4.2.14 comes from the use of \( q : k[x] \to k[x] \) (see also the proof of Theorem 4.2.22 for an example of this fact).

Let us consider the following situation:

\[
\begin{array}{c}
k[x] \xrightarrow{\varphi} k[y] \\
\downarrow q \quad \downarrow \varphi \\
k[x]
\end{array}
\]

where \( \varphi \) is injective, \( r_1(\varphi) = n - 1 \) and \( q \) defined by \( q(x_1) = x_1x_2 \) and \( q(x_i) = x_i \) for \( i \neq 1 \). It may happen that \( \varphi \) it is not injective. In this case let \( z \) be a generator of \( \text{Ker}(\varphi) \). Since \( \varphi \) is injective and \( \varphi \) is not injective, we can prove that \( z \) is not algebraic over \( q(k[x]) = k[x_1x_2, x_2, ..., x_n] \) (see Part 4.2 [Ron09b]). Moreover the fact that the Chevalley function of \( \varphi \) is not bounded by an affine function comes from the following transcendence property of \( z \) (see Lemma 4.7 [Ron09b]):

\textbf{Lemma 4.2.19.} — There exists a decreasing function \( \alpha : \mathbb{R}^+ \to \mathbb{R}^+ \) such that

\[
\left| \frac{f}{g} - z \right| \geq \alpha(|g|) \quad \forall f \in k[x_1x_2, x_2, ..., x_n], \ g \in k[x_1x_2, x_2, ..., x_n][x_1],
\]

and if \( \alpha \) is the greatest function satisfying the above inequality, then \( \frac{\ln(\alpha(u))}{\ln(\alpha)} \to 0 \) as \( u \) goes to 0.

\textbf{4.2.6. Two "good" examples.} — Now we can give two important examples of regular morphisms: these are the analytic morphisms defined by algebraic power series (Remark 4.2.18 already gave a flavor of this) and the morphisms whose source is a domain of dimension less or equal to two.

\textit{4.2.6.1. Case of the algebraic morphisms.} — Let us denote by \( k\langle x \rangle := k\langle x_1, ..., x_n \rangle \) the subring of \( k[[x_1, ..., x_n]] \) of all power series which are algebraic over \( k[x_1, ..., x_n] \). We clearly have \( k\langle x \rangle \subset k\{x\} \) for any \( n \in \mathbb{N} \) when \( k \) is a valued field. We call Henselian \( k \)-algebra any \( k \)-algebra which is isomorphic to a quotient \( k\langle x \rangle/I \) where \( n \in \mathbb{N} \) and \( I \) is an ideal of \( k\langle x \rangle \). If \( \varphi : k\langle x \rangle/I \to k\langle y \rangle/J \) is a morphism of Henselian algebras, we denote by \( \tilde{\varphi} : k\{x\}/Jk\{x\} \to k\{y\}/Jk\{y\} \) the induced analytic morphism. If \( \varphi : A \to B \) is a morphism of Henselian algebras we define \( r_4(\varphi) := \dim(A/\text{Ker}(\varphi)) \). We also denote \( r_i(\varphi) := r_i(\tilde{\varphi}) \) for \( 1 \leq i \leq 3 \). We obviously have \( r_1(\varphi) \leq r_2(\varphi) \leq r_3(\varphi) \leq r_4(\varphi) \).

We remind that any polynomial \( P \) with coefficients in \( k\langle x, y \rangle \) satisfies the
nested Artin approximation property (cf. Theorem 3.5.8). This means that any solution \((f_i, g_j) \in k[[x]]^r \times k[x, y]^s\) of \(P(F_i, G_j) = 0\) may be approximated by solutions in \(k[[x]]^r \times k(x, y)^s\). By using Proposition 4.2.7 that is valid in the Henselian case, we see that any morphism \(\varphi : k(x) \rightarrow k(y)\) satisfies \(\widehat{\varphi}(k[[x]]) \cap k(y) = \varphi(k(x))\). In fact we can show more than that by using only Theorem 4.2.16 which is valid for rings of algebraic power series and not Theorem 3.5.8.

**Theorem 4.2.20.** — [Ron09b] Let \(\varphi : A \rightarrow k\langle y \rangle\) be a morphism of Henselian \(k\)-algebras. Then \(r_1(\varphi) = r_4(\varphi)\).

**Proof.** — We can easily reduce the problem to the case \(\varphi\) is injective and \(A = k\langle x \rangle\).

Then we apply Theorem 4.2.16 which is still valid if we replace \(k\{x\}\) and \(k\{y\}\) by \(k\langle x \rangle\) and \(k\langle y \rangle\) (since the rings of algebraic power series form a Weierstrass system). Thus we have a commutative diagram:

\[
\begin{array}{ccc}
k(x) & \xrightarrow{\varphi} & k(y) \\
\downarrow{\sigma_1} & & \downarrow{\sigma_2} \\
k(x) & \xrightarrow{\pi} & k(y)
\end{array}
\]

We see that \(\sigma_2 \circ \varphi\) is injective, thus \(\pi \circ \sigma_1\) is injective. Thus it is enough to show the following: if \(\overline{\varphi} : k\langle x \rangle \rightarrow k\langle y \rangle\) is a non injective morphism satisfying \(\varphi = \overline{\varphi} \circ \sigma\) where \(\sigma\) is one of the three types of morphisms defined in Theorem 4.2.16 i), then \(\varphi\) is not injective. This is straightforward to check this for any of these three morphisms. \(\square\)

Next result has been proven by Tougeron [To76], Becker [Be77a] and Milman [Mi78a] in the case of complex analytic algebras which are the quotient of the ring of convergent power series by an ideal generated by polynomials. We extended this result for quotients of convergent power series rings by ideals generated by algebraic power series in any characteristic.

**Corollary 4.2.21.** — [To76] [Be77a] [Mi78a] [Ron06b] Let \(\varphi : k\{x\}/Jk\{x\} \rightarrow k\{y\}/Jk\{y\}\) be a morphism of analytic \(k\)-algebras where \(I\) is an ideal of \(k\langle x \rangle\), \(J\) a prime ideal of \(k\langle y \rangle\) and such that \(\varphi(x_i) \in k\langle y \rangle/J\) for any \(1 \leq i \leq n\). Let us assume that \(\text{char}(k) = 0\) or \(J = (0)\). Then \(r_1(\varphi) = r_3(\varphi)\).

**Proof.** — Let \(\pi : k\{x\} \rightarrow k\{x\}/I\) be the quotient morphism. We can replace \(\varphi\) by the composed morphism \(\varphi \circ \pi : k\{x\} \rightarrow k\langle y \rangle/Jk\{y\}\) and we still denote this morphism by \(\varphi\). We denote by \(\varphi^h : k\langle x \rangle \rightarrow k\langle y \rangle/J\) the associated
morphism of Henselian algebras and we see that \( r_1(\varphi^h) = r_4(\varphi^h) \) by the previous theorem if \( J = (0) \). If \( \text{char}(\mathbb{k}) = 0 \) we can use resolution of singularities as explained in Remark 4.2.15. Hence \( r_3(\varphi) = r_3(\varphi^h) = r_1(\varphi^h) = r_1(\varphi) \).

4.2.6.2. The dimension 2 case. —

**Theorem 4.2.22.** — [Ron09b] Let \( \varphi : A \to B \) be a morphism of analytic \( \mathbb{k} \)-algebras where \( A \) is an integral domain of dimension 2 and \( B \) is regular. Then \( \varphi \) is injective if and only if \( r_1(\varphi) = 2 \). This result remains valid if \( A = \mathbb{k}\langle x \rangle \) and \( B = \mathbb{k}\langle y \rangle \) where the family \( (\mathbb{k}\langle x_1, \ldots, x_n \rangle)_{n \in \mathbb{N}} \) is a Weierstrass system over any field \( \mathbb{k} \).

**Sketch of the proof.** — If \( r_1(\varphi) = 2 \) then \( r_3(\varphi) = 2 \) and \( \varphi \) is injective. The non trivial part of the theorem is the reverse implication.

First of all we may reduce to the case \( A = \mathbb{k}\{x_1, x_2\} \) and \( B = \mathbb{k}\{y\} \). Then we apply Theorem 4.2.16 for the morphism \( \varphi \). Thus we have the following commutative diagram:

\[
\begin{array}{ccc}
\mathbb{k}\{x_1, x_2\} & \xrightarrow{\varphi} & \mathbb{k}\{y\} \\
\downarrow{\sigma_1} & & \downarrow{\sigma_2} \\
\mathbb{k}\{x_1, x_2\} & \xrightarrow{\overline{\varphi}} & \mathbb{k}\{y\}
\end{array}
\]

If \( \sigma_1 \) is uniquely a composition of \( \mathbb{k} \)-automorphisms and of morphisms \( \chi_d \) (\( d \in \mathbb{N}^\ast \) being a prime number) defined by \( \chi_d(x_1) = x_1^d \), and \( \chi_d(x_i) = x_i \) \( \forall i \neq 1 \), and if \( \varphi \) is injective then it is not too difficult to see that \( \overline{\varphi} \) is still injective and thus \( r_1(\varphi) = 2 \).

Now the idea is to analyze the proof of Theorem 4.2.16 in order to see that, if \( \varphi \) is injective, we can construct a commutative diagram similar to the previous one but where \( \sigma_1 \) is uniquely a composition of \( \mathbb{k} \)-automorphisms and of morphism \( \chi_d \) and where \( \overline{\varphi} \) is defined by

\[
\overline{\varphi}(x_1) = y_1^a y_2^b u \\
\overline{\varphi}(x_2) = y_1^c y_2^d v
\]

such that the rank of the matrix \( \begin{pmatrix} a & b \\ c & d \end{pmatrix} \) is equal to 2 and \( u \) and \( v \) are units.

Hence \( r_1(\overline{\varphi}) = 2 \), so \( r_1(\varphi) = 2 \).

One corollary of this theorem is the fact that such an injective morphism \( \varphi : A \to B \) where \( \dim(A) = 2 \) and \( B \) is regular satisfies \( \varphi(A) = \hat{\varphi}(A) \cap B \). This corollary (in the case of morphisms between convergent power series rings) has been proven before in [AvdP70].
Corollary 4.2.23. — Let $\varphi : A \to B$ be a morphism of analytic $k$-algebras where $A$ is an integral domain of dimension $\leq 2$ and $B$ is regular. Then $r_1(\varphi) = r_2(\varphi) = r_3(\varphi)$.

Proof. — If $\dim(A) = 0$, this is trivial.
If $\dim(A) = 1$ and $r_1(\varphi) = 0$, then $\Ker(\varphi) = m_A$ thus $r_3(\varphi) = 0$. If $r_1(\varphi) = 1$ then $r_3(\varphi) = 1$.
If $\dim(A) = 2$ and $r_1(\varphi) = 0$, then $\Ker(\varphi) = m_A$ so $r_3(\varphi) = 0$. If $r_3(\varphi) = 2$ then $\varphi$ is injective and $r_1(\varphi) = 2$ by the previous theorem. Finally if $r_1(\varphi) = 1$ then $r_3(\varphi) = 1$ since $r_3(\varphi) = 2$ implies $r_1(\varphi) = 2$.

Remark 4.2.24. — Remark 4.2.15 remains valid here and we may assume that $B$ is just an integral domain when $\text{char}(k) = 0$ in both results.

4.3. Algebraic closure of the field of power series in several variables in characteristic zero

This work has been motivated by understanding the proof of Gabrielov Theorem given by J.-C. Tougeron [To90] and by a question of A. Parusiński.

The last part of this thesis concerns the problem of "describing" the set of solutions of polynomial equations with coefficients in the formal power series ring $k[[x_1, \ldots, x_n]]$. Here we are interested by the set of roots of one polynomial in one variable with coefficients in $k[[x_1, \ldots, x_n]]$. Let us begin by surveying the case of polynomials with coefficients in the field of power series in one variable. In the whole part we will assume that the characteristic of the base field $k$ is equal to zero.

When $k$ is an algebraically closed field of characteristic zero, we can always express the roots of a polynomial with coefficients in the field of power series over $k$, denoted by $k((t))$, as formal Laurent series in $t^\frac{1}{k}$ for some positive integer $k$. This result was known by Newton (at least formally see [BK86] p. 372) and had been rediscovered by Puiseux in the complex analytic case [Pu50], [Pu51] (see [BK86] or [Cu04] for a presentation of this result). A modern way to reformulate this fact is to say that an algebraic closure of $k((t))$ is the field of Puiseux power series $\mathcal{P}$ defined in the following way:

\[ \mathcal{P} := \bigcup_{k \in \mathbb{N}} k \left( t^\frac{1}{k} \right). \]

The proof of this result, called the Newton-Puiseux method, consists essentially in constructing the roots of a polynomial $P(Z) \in k[t][Z]$ by successive approximations (see [BK86] or [Cu04] - see [BiMi90] for a slightly different
method). These approximations converge since \( \mathbb{k}\left(\left(\frac{1}{t^k}\right)\right) \) is a complete field with respect to the Krull topology (see Remark 3.2.3 for a definition of this topology).

The result is quite impressive since it gives also a description of the Galois group of \( \mathbb{k}(\left(t^\frac{1}{k}\right)) \rightarrow \mathbb{P} \). Indeed this one is generated by the multiplication of the \( k \)-th roots of unity with \( t^\frac{1}{k} \) for any positive integer \( k \). In particular the conjugates of any convergent power series in \( \mathbb{C}\{t^\frac{1}{k}\} \) are also in \( \mathbb{C}\{t^\frac{1}{k}\} \). This has some important corollaries as the fact that if an irreducible polynomial with coefficients in \( \mathbb{C}[t] \) has a root which is a convergent power series in \( t^\frac{1}{k} \) the the others roots are also convergent power series.

When \( \mathbb{k} \) is a characteristic zero field (not necessarily algebraically closed), we can prove in the same way that an algebraic closure of \( \mathbb{k}(\left(t^\frac{1}{k}\right)) \) is

\[
\mathbb{P} := \bigcup_{k'} \bigcup_{k \in \mathbb{N}} \mathbb{k}'\left(\left(\frac{1}{t^k}\right)\right).
\]

where the first union runs over all finite field extensions \( \mathbb{k} \rightarrow \mathbb{k}' \).

It is tempting to find such a similar result for the algebraic closure of the field of power series in \( n \) variables, \( \mathbb{k}(\left(x_1, \ldots, x_n\right)) \), for \( n \geq 2 \). But it appears easily that the algebraic closure of this field admits a really more complicated description and considering only power series depending on \( x_1^\frac{1}{k}, \ldots, x_n^\frac{1}{k} \) is not sufficient. For instance it is easy to see that a root square of \( x_1 + x_2 \) can not be expressed as such a power series.

Nevertheless there exist positive results in some specific cases, the more famous one being the Abhyankar-Jung theorem:

4.3.1. Abhyankar-Jung Theorem. —

**Theorem 4.3.1 (Abhyankar-Jung Theorem). —** If \( \mathbb{k} \) is a field of characteristic zero, then any monic polynomial with coefficients in \( \mathbb{k}[x_1, \ldots, x_n] \), whose discriminant has the form \( ux_1^{\alpha_1} \ldots x_n^{\alpha_n} \) where \( u \in \mathbb{k}[x_1, \ldots, x_n] \) is a unit and \( \alpha_1, \ldots, \alpha_n \in \mathbb{Z}_{\geq 0} \), has its roots in \( \mathbb{k}'[x_1^\frac{1}{k}, \ldots, x_n^\frac{1}{k}] \) where \( k \in \mathbb{N}^* \) and \( \mathbb{k} \rightarrow \mathbb{k}' \) is a finite field extension.

This result can be seen as a generalization of Newton-Puiseux Theorem since any polynomial with coefficients in the ring of powers series in one variable satisfies the hypothesis of Abhyankar-Jung Theorem. This result has first been proven by Jung in the complex analytic case [Ju08] then by Abhyankar in the general case. A polynomial satisfying the hypothesis of this theorem is called a quasi-ordinary polynomial.
Jung’s proof is very elementary but uses a topological argument. Here follows a sketch of this proof:

**Sketch of the proof of the Abhyankar-Jung Theorem in the analytic case**

Let $U$ be the following polydisc:

$$
U = \prod_{i=1}^{n} D_{\varepsilon} = \{ x \in \mathbb{C}^n ; |x_i| < \varepsilon, i = 1, \ldots, n \}
$$

with $\varepsilon$ small enough such that the coefficients of $P(Z)$ are analytic in a neighborhood of the closure of $U$. By hypothesis the projection of $\{(x_1, \ldots, x_n, Z) \in U \times \mathbb{C} ; P(x, Z) = 0 \}$ over $U$ is a finite covering which is ramified over the hyperplanes of coordinates. By the lifting homotopical criterion there exists a positive integer $q$ such that the map

$$
x(y) = (y_1^q, \ldots, y_n^q) : U_1 \to U, \quad \text{where } U_1 = \prod_{i=1}^{n} D_{\varepsilon^1/q},
$$

induces a trivial covering over $U_1^* = \prod_{i=1}^{n} D_{\varepsilon^1/q}^*$. This is equivalent to say that the roots of $P(x(y), Z)$ are analytic on $U_1^*$. Moreover the polynomial $P(Z)$ being monic its roots are bounded in a neighborhood of the origin, thus they may be extended to an analytic function on $U_1$.

Abhyankar’s proof is purely algebraic and it is not easy to understand. Thirty years ago Luengo published a paper presenting a new proof of Abhyankar-Jung Theorem in the general case which was more elementary. This proof used a property of the Newton polyhedron satisfied by quasi-ordinary polynomials, and the fact that quasi-ordinary polynomials satisfy this property is equivalent to the Abhyankar-Jung Theorem. Unfortunately it appeared that there was a serious gap in the proof. This motivated Kiyek and Vicente to give a new proof, purely algebraic, of the Abhyankar-Jung Theorem based on the theory of ramified morphisms between local rings [KiVi04].

With Adam Parusiński we gave an elementary proof of the result announced by Luengo. In order to present this result we will first give a definition introduced by Hironaka:

**Definition 4.3.2** — [Hir74] Let $P(Z) \in k[[x_1, \ldots, x_n]][Z]$ be a monic polynomial of degree $d$ and let $NP$ denote its Newton polyhedron. Let us write

$$
P(Z) = \sum_{(i_1, \ldots, i_{n+1}) \in \mathbb{N}_{n+1}} P_{i_1, \ldots, i_{n+1}} x_1^{i_1} \ldots x_n^{i_n} Z^{i_{n+1}}.
$$

This polynomial is called $\nu$-quasi-ordinary if there is a point $R_1$ of the Newton polyhedron $NP$, $R_1 \neq R_0 = (0, \ldots, 0, d)$, such that if $R'_1$ denotes the projection
of $R_1$ onto $\mathbb{R}^n \times \{0\}$ from $R_0$, and $S = |R_0, R'_1|$ is the segment joining $R_0$ and $R'_1$, then

$$NP \subset |S| := \bigcup_{s \in S} (s + \mathbb{R}^{n+1}_{\geq 0})$$

and

$$P_S = \sum_{(i_1, \ldots, i_{n+1}) \in S} P_{i_1, \ldots, i_{n+1}} x_1^{i_1} \cdots x_n^{i_n} Z^{i_{n+1}}$$

is not a power of degree one polynomial in $Z$.

Last condition is automatically satisfied if the coefficient of $Z^{d-1}$ in $P(Z)$ is zero since $k$ is a field of characteristic zero.

Here is a picture of the Newton polyhedron of a $\nu$-quasi-ordinary polynomial with $n = 2$ (thick lines represent the edges of the Newton-Polyhedron):

\[ \text{Diagram of Newton polyhedron} \]

Then the next result is equivalent to the Abhyankar-Jung Theorem:

**Theorem 4.3.3.** — [PaRo12] Let $P(Z) \in k[[x_1, \ldots, x_n]][Z]$ be a monic polynomial of degree $d$. If $P(Z)$ is quasi-ordinary and if the coefficient of $Z^{d-1}$ in $P(Z)$ is zero then $P(Z)$ is $\nu$-quasi-ordinary.

**Sketch of the proof.** — In [PaRo12] we give two different proofs of this result. Both are based on the complex analytic case of Abhyankar-Jung Theorem. We present here one of these proofs which uses Artin approximation Theorem. We write $P(Z) = Z^d + a_2 Z^{d-2} + \cdots + a_d$.

- **First step:** By using Jung’s proof sketched before we see that the Abhyankar-Jung theorem is true for polynomials whose coefficients are convergent power series over $\mathbb{C}$.
- **Second step:** Then we show that if the discriminant of a polynomial $P(Z)$ has normal crossing and if its roots are Puiseux series in several variables and
$a_1 = 0$, then $P(Z)$ is $\nu$-quasi-ordinary. In fact these three facts imply that the set of non-zero exponents appearing in the Taylor expansion of the roots of $P(Z)$ is totally ordered in $\mathbb{Q}^n$. This is not very difficult and is based on a short combinatorial study of these exponents. Thus these roots are equal to a monomial times a unit. For any $k$, $a_k$ is a homogeneous function of degree $k$ in the roots of $P(Z)$, thus the ideal $(a_1^d, a_2^d, \ldots, a_d^d)$ is a monomial ideal generated by one these $d$ monomials. This is equivalent to say that $NP \subset |S| := \bigcup_{s \in S} (s + \mathbb{R}_{\geq 0}^{n+1})$. Thus Theorem 4.3.3 is proven for polynomials with complex analytic coefficients.

- **Last step**: Now we prove Theorem 4.3.3 in the general case. Let us consider a monic polynomial $P(Z) \in \mathbb{k}[x_1, \ldots, x_n][Z]$ whose discriminant is normal crossing (i.e. of the form $u x^\alpha$ where $u$ is a unit and $\alpha \in \mathbb{N}^n$) and where $\mathbb{k}$ is a characteristic zero field. Since the coefficients of $P(Z)$ depend only on a countable number of elements of $\mathbb{k}$ we may assume that $\mathbb{Q} \longrightarrow \mathbb{k}$ is an extension whose transcendence degree is at most countable. Such a field extension embeds in $\mathbb{C}$. Thus we may assume that the coefficients of $P(Z)$ are power series over $\mathbb{C}$ since this does not change the shape of its Newton polyhedron.

The discriminant of $P(Z)$ is a polynomial $\Delta = \Delta(a_2, \ldots, a_d)$ depending on the coefficients $a_2, \ldots, a_d$. Now let us define the following polynomial:

$$Q(A_2, \ldots, A_d, U) := \Delta(A_2, \ldots, A_d) - U x^\alpha \in \mathbb{C}\{x\}\{A_2, \ldots, A_d, U\}.$$ 

Then $Q(a_2, \ldots, a_d, u) = 0$ and by Artin approximation Theorem we may find, for any integer $c$, convergent power series $a_{2,c}, \ldots, a_{d,c}$, $u_c$ solutions of $Q = 0$ and equal to $a_2, \ldots, a_d, u$ up to order $c$. In particular the polynomial $P_c(Z) := Z^d + a_{2,c}Z^{d-2} + \cdots + a_{d,c} \in \mathbb{C}\{x\}\{Z\}$ is quasi-ordinary. Thus $P_c(Z)$ is $\nu$-quasi-ordinary by the previous step. Since the coefficients of $P_c(Z)$ coincide with those of $P(Z)$ up to order $c$, the Newton polyhedron of $P_c(Z)$ is included in the Newton polyhedron of $P_c(Z)$ modulo high terms. But the Newton polyhedron of $P_c(Z)$ is included in $|S|$ where $S$ does not depends on $c$ if $c$ is large enough (if $c$ is larger than the size of the vertices of the Newton polyhedron of $P(Z)$). Thus, at the limit, the Newton polyhedron of $P(Z)$ is included in $|S|$. Hence $P(Z)$ is $\nu$-quasi-ordinary.

We remark that the second step proves that the Abhyankar-Jung Theorem implies Theorem 4.3.3. This was known before. The fact that Theorem 4.3.3 implies the Abhyankar-Jung Theorem comes from the fact that if the coefficient of $Z^{d-1}$ of the quasi-ordinary polynomial $P(Z)$ is zero (we can always assume this after Tschirnhaus transform $Z \mapsto Z - \frac{a_d}{d!}$) then its Newton polyhedron "begins" with one face of dimension 1 et this allows us to repeat the classical Newton-Puiseux method for polynomials with coefficients in the field of power series in one variables. Thus both theorems are equivalent.
In particular this allows us to prove the following result that says that the roots of a quasi-ordinary polynomial whose coefficients are in a "good" subring of power series are still in this subring (after replacing the $x_i$’s by some of their powers):

**Theorem 4.3.4.** — [PaRo12] Let $k$ be an algebraically closed field of characteristic zero. Let us consider, for any $n \in \mathbb{N}$, a subring $k\{\{x\}\}$ of $k[[x]]$ containing the ring of polynomials, satisfying the Implicit function Theorem, stable by division by the $x_i$’s and stable par composition by powers of the $x_i$’s. Then any monic polynomial with coefficients in $k\{\{x_1, ..., x_n\}\}$ whose discriminant is equal to $ux_1^{\alpha_1}...x_n^{\alpha_n}$, where $u \in k\{\{x_1, ..., x_n\}\}$ is a unit and $\alpha_1, ..., \alpha_n \in \mathbb{N}$, has its roots in $k\{\{x_1^k, ..., x_n^k\}\}$ for some $k \in \mathbb{N}^*$.

This result is in particular valid for the rings of germs of quasi-analytic function that do not satisfy the Weierstrass preparation theorem.

Let us mention also that we can use the strong Artin approximation Theorem (see Corollary 3.3.16) instead of the classical Artin approximation Theorem in the last step. In this case we obtain the following result (saying that if $P(Z)$ is close to be quasi-ordinary then it is close to be $\nu$-quasi-ordinary):

**Theorem 4.3.5.** — [PaRo12] Let $d \in \mathbb{N}$ and $\alpha \in \mathbb{N}^n$. Then there exists a function $\beta : \mathbb{N} \rightarrow \mathbb{N}$ satisfying the following property: For any integer $c$ and any monic polynomial $P(Z) \in k[[x_1, ..., x_n]()[Z]$ of degree $d$ in $Z$ whose discriminant is equal to $x_1^{\alpha_1}...x_n^{\alpha_n}$ times a unit modulo $(x)^{\beta(c)}$, there exists a compact face of dimension 1 of the Newton polyhedron of $P(Z)$ containing $(0, ..., 0, d)$, denoted by $S$, such that

$$NP \subset |S| + \{j \in \mathbb{N}^n/j_1 + \cdots + j_n \geq c\}.$$ 

Finally let us mention that this method allows us to prove the Abhyankar-Jung Theorem for polynomials with coefficients in $k[[x_1, ..., x_n]]$ where $k$ is a characteristic zero field and $I$ is a binomial ideal (see Theorem 6.2 [PaRo12]). This result generalizes a theorem proven by Gonzálež Pérez in the case of polynomials with coefficients in the ring of germs of holomorphic functions in a point of toric variety [Go00].

**4.3.2. Newton-Puiseux method for Abhyankar valuations and generalization of the Abhyankar-Jung Theorem.** — The second work done in relations with the description of the roots of polynomials with coefficients in the ring of power series over a field of characteristic zero is the study of the Newton-Puiseux method with respect to a rank one Abhyankar valuation.

The first natural idea to find the roots of a polynomial with coefficients in
the ring of power series in $n$ variables over a field of characteristic zero involves the use of Newton-Puiseux theorem $n$ times (i.e. the formula \([14]\) for the algebraic closure of $k((t))$). For example in the case $n = 2$, this means that the algebraic closure of $k((x_1, x_2))$ is included in

$$L := \bigcup_{k_2 \in \mathbb{N}} \bigcup_{k_1 \in \mathbb{N}} k \left( \left( \frac{1}{x_1^k} \right), \left( \frac{1}{x_2^k} \right) \right).$$

But this field, which is algebraically closed, is very much larger than the algebraic closure of $k((x_1, x_2))$ (see \([Sa10]\) for some thoughts about this). Moreover the action of the $k_1$-th and $k_2$-th roots of unity are not sufficient to generate the Galois group of the algebraic closure since there exist elements of $k((x_1, x_2))$ which are algebraic over $k((x_1, x_2))$ but are not in $k((x_1, x_2))$. For instance consider

$$x_1 \sqrt{1 + \frac{x_2}{x_1}} = \sum_{i \geq 0} \left( \frac{1}{i} \right) \frac{1}{x_1^{i-1} x_2} \in \mathbb{Q}((x_1))((x_2)) \setminus \mathbb{Q}((x_1, x_2)).$$

Nevertheless a deeper analysis of the Newton-Puiseux method leads to the fact that it is enough to consider the field of fractions of Puiseux power series whose support is included in a rational strictly convex cone:

**Theorem 4.3.6.** — \([McD95]\) Let $P(Z) \in k((x_1, \ldots, x_n))[Z]$. Then there exist a rational strictly convex cone $\sigma$ containing $\mathbb{R}_{\geq 0}^n$ and $k \in \mathbb{N}^*$ such that the roots of $P(Z)$ are in the fraction field of

$$\left\{ f = \sum_{(l_1, \ldots, l_n) \in \mathbb{Z}^2} a_{l_1, \ldots, l_n} \frac{t_1^{l_1}}{x_1^{l_1}} \ldots \frac{t_n^{l_n}}{x_n^{l_n}} / \text{Supp}(f) \subset \sigma \right\}.$$

We can also find a proof of this result and of some strengthened versions of it in \([Go00]\), \([Aro04]\), \([ArIl09]\), \([SV11]\). But once more, for any rational strictly convex cone of $\mathbb{R}^2$, denoted by $\sigma$, $\mathbb{R}_{\geq 0}^2 \subset \sigma$, there exist elements whose support is in $\sigma$ but that are not algebraic over $k((x_1, \ldots, x_n))$. On the other hand if a power series with support in $\sigma$ is algebraic over $k((x_1, \ldots, x_n))$ it is not clear what are its conjugates.

Hence we see that two problems emerge:

- Characterize the power series with support in "large" cones that are algebraic over $k((x_1, \ldots, x_n))$ (or at least give necessary conditions in order to insure that such a series is algebraic over $k((x_1, \ldots, x_n))$).

- Find a description of the Galois group of a polynomial $P(Z)$ with coefficients in $k[[x_1, \ldots, x_n]]$, or at least relate properties of one root of $P(Z)$ with
properties of the others.

The work done in [Ron13b] and presented here is strongly inspired by the article [To90] where a similar study is done for the valuation ord when $k = \mathbb{C}$ with analytic methods. Most of the results presented here have been stated by Tougeron in the case of the valuation ord. Our work has essentially two parts. First of all we study the Newton-Puiseux method for rank one Abhyankar val-

We denote by $k[[x]]$ the ring of formal power series in $n$ variables over a characteristic zero field $k$. We consider a rank one valuation $\nu$ which is non-negative on $k[[x]]$ and centered at the maximal ideal of $k[[x]]$. Such valuation is called an Abhyankar valuation if the Abhyankar inequality is an equality for it. This is equivalent to say that it is a monomial valuation after a sequence of blowing-ups. We denote by $V_\nu$ the valuation ring of $\nu$ and by $\hat{V}_\nu$ its completion for the Krull topology. We denote by $\mathbb{K}_n$ the faction field of $k[[x]]$ and by $\hat{\mathbb{K}}_\nu$ that of $\hat{V}_\nu$. We set $\text{Gr}_\nu V_\nu := \bigoplus_{i \in \mathbb{R}^+} \frac{p_{\nu,i}}{p_{\nu,i}^+}$, where $p_{\nu,i} := \{ f \in V_\nu / \nu(f) \geq i \}$ and $p_{\nu,i}^+ := \{ f \in V_\nu / \nu(f) > i \}$, the associated graded ring of $V_\nu$. Since $\nu$ is an Abhyankar valuation we can show that we have an isomorphism

$$\hat{V}_\nu \simeq \text{Gr}_\nu V_\nu.$$ In particular we can see $k[[x]]$ and $V_\nu$ as subrings of $\text{Gr}_\nu V_\nu$.

In order to apply the Newton-Puiseux method in $k[[x]]$, consisting essentially to construct a root of a monic polynomial of $k[[x]][Z]$ by successive approximations, we need to work in the associated graded ring of $V_\nu$ and define what will be the equivalent of the fractional powers of $t$ when we use the Newton-Puiseux method in $k[[t]]$ (the ring of power series in one variable). This is the motivation of the following definition:

**Definition 4.3.7.** [Ron13b] Let us fix $d \in \mathbb{R}^+$. A homogeneous element (of degree $d$) with respect to the valuation $\nu$, is an element $\gamma$ belonging to a finite extension of $V_\nu$ whose minimal polynomial has the form

$$Z^q + g_1 Z^{q-1} + \cdots + g_q$$

where $g_k \in p_{\nu,dk}$ for $1 \leq k \leq q$. If $g_k$ is the image of an element of $k[[x]]$, for all $k$, we say that $\gamma$ is an integral homogeneous element with respect to $\nu$.

In this case the valuation $\nu$ extends uniquely to $\text{Gr}_\nu V_\nu[\gamma]$ by defining $\nu(\gamma) = d$. 


**Example 4.3.8.** — If $\nu$ is a monomial valuation whose weights $\alpha_1, \ldots, \alpha_n$ are $\mathbb{Q}$-linearly independent, then the homogeneous elements with respect to $\nu$ are the monomials of the form $x_1^{\beta_1} \cdots x_n^{\beta_n}$ where $(\alpha, \beta) \geq 0$ and $\beta_i \in \mathbb{Q}$.

Integral homogeneous elements are those for which the $\beta_i$’s are non-negative.

**Example 4.3.9.** — If $c$ is algebraic over $k$ then $c$ is an integral homogeneous element of degree 0 with respect to any Abhyankar valuation.

Now we can apply Newton-Puiseux method and we obtain the following result:

**Theorem 4.3.10.** — [Ron13b] Let $k$ be a characteristic zero field and let $\nu$ be an Abhyankar valuation. Set $N := \text{dim}_Q \Gamma \otimes \mathbb{Z} \mathbb{Q}$. For any $P(Z) \in k((x))[Z]$ there exist integral homogeneous elements $\gamma_1, \ldots, \gamma_N$ with respect to $\nu$ such that the roots of $P(Z)$ are in $\hat{K}_\nu[\gamma_1, \ldots, \gamma_N]$.

The fact that we need only $N$ homogeneous elements comes from the Primitive Element Theorem.

This result asserts that the inductive limit of the fields $\hat{K}_\nu[\gamma_1, \ldots, \gamma_N]$ when $\gamma_1, \ldots, \gamma_N$ run over all integral homogeneous elements (we denote this limit by $\mathbb{K}_\nu$) contains an algebraic closure of $k((x))$ (in fact it is an algebraically closed field). Thus we have two field extensions:

\[ k((x)) \longrightarrow \hat{K}_\nu \longrightarrow \mathbb{K}_\nu \]

and the Galois group of the second extension has a quite simple description since it acts only on homogeneous elements. Thus it is very natural to study irreducible monic polynomials of $k[[x]][Z]$ that remain irreducible in $\hat{V}_\nu[Z]$ since their Galois group will act only on homogeneous elements.

Let us mention that we can also prove the following result that will be very useful in the sequel (for a polynomial $R(Z) \in \hat{V}_\nu[Z]$, we say that $\nu(R(Z)) \geq r$ if all the coefficients $a_i$ of $R(Z)$ satisfy $\nu(a_i) \geq r$):

**Proposition 4.3.11.** — [To90, Ron13b] Let $P(Z) \in \hat{V}_\nu[Z]$ be a monic polynomial of degree $d$ without no multiple factor. Let us write $P(Z) = P_1(Z) \cdots P_r(Z)$ where the $P_i(Z) \in \hat{V}_\nu[Z]$ are irreducible monic polynomials. Let $Q(Z) \in \hat{V}_\nu[Z]$ be another monic polynomial of degree $d$. Let $z_1, \ldots, z_d$ be the roots of $P(Z)$. If

\[ \nu(P(Z) - Q(Z)) > d \max_{i \neq j}\{\nu(z_i - z_j)\} \]
then we may factor $Q(Z) = Q_1(Z)...Q_r(Z)$ where the $Q_i(Z) \in \hat{V}_\nu[Z]$ are irreducible monic polynomials and

$$\nu(Q_i(Z) - P_i(Z)) \geq \frac{\nu(P(Z) - Q(Z))}{d}.$$ 

**Remark 4.3.12.** — The hypothesis is satisfied if $\nu(P(Z) - Q(Z)) > \frac{d}{2} \nu(\Delta_P)$ where $\Delta_P$ is the discriminant of $P(Z)$.

Then we study the particular case of monomial valuations whose weights are positive integers. We can prove that in the previous construction we can replace $\hat{K}_\nu$ by a smaller field. First let us give a definition: we fix $\alpha \in \mathbb{R}^n_{>0}$ and we denote by $\nu_\alpha$ the monomial valuation defined by $\nu_\alpha(x_i) = \alpha_i$ for all $i$. A $(\alpha)$-homogeneous polynomial is a weighted homogeneous polynomial for the weights $\alpha_1,...,\alpha_n$. Then we can define the following valuation ring:

$$V_\alpha := \left\{ A \in \hat{V}_{\nu_\alpha} \mid \exists \Lambda \text{ a finitely generated sub-semigroup of } \mathbb{R}_{\geq 0}, \right.$$

$$\exists \theta \in \mathbb{k}[x] \text{ (} \alpha \text{-homogeneous), } \forall i \in \Lambda \exists a_i \in \mathbb{k}[x] \text{ (} \alpha \text{-homogeneous), }$$

$$\exists a \geq 0, b \in \mathbb{R} \forall i \in \Lambda \exists m(i) \in \mathbb{N} \text{ s.t. } m(i) \leq ai + b, \nu_\alpha \left( \frac{a_i}{\theta^m(i)} \right) = i \text{ and } A = \sum_{i \in \Lambda} a_i \theta^m(i) \right\}.$$ 

We denote by $K_\alpha$ its fraction field. Thus we have the following result whose proof is inspired by a result of Gabrielov [Gab73] and is based on the Implicit function Theorem of Tougeron (cf. Theorem 3.3.2).

**Theorem 4.3.13.** — [To90][Ron13b] Let $k$ be a characteristic zero field and let $\nu_\alpha$ be a monomial valuation. We set $N := \dim_Q \Gamma \otimes_{\mathbb{Z}} Q$. For any $P(Z) \in k((x))[Z]$ there exist integral homogeneous elements $\gamma_1,...,\gamma_N$ with respect to $\nu$ such that the roots of $P(Z)$ are in $K_\alpha[\gamma_1,...,\gamma_N]$.

In the case the $\alpha_i$'s are $\mathbb{Q}$-linearly independent this statement is exactly Theorem 4.3.6 and the cone $\sigma$ satisfies $\langle \beta, \alpha \rangle > 0$ for any $\beta \in \sigma$, $\beta \neq 0$ (this comes essentially from Example 4.3.8).

In the case $k = \mathbb{C}$, $\alpha \in \mathbb{N}^n$ and the coefficients of $P(Z)$ are in $\mathbb{C}\{x\}$, we can replace $V_\alpha$ by the following valuation ring $V_\alpha^{\mathbb{C}\{x\}}$ (cf. Example 6.13 [Ron13b]):

$$V_\alpha^{\mathbb{C}\{x\}} := \left\{ A = \sum_{i \in \Lambda} a_i \theta^m(i) \in V_\alpha \mid \exists C, r > 0 \text{ t.q. } |a_i(\xi)| \leq C r^i ||\xi||_{\nu_\alpha}(a_i) \forall \xi \in \mathbb{C}^n \right\}.$$
where $||\xi||_\alpha := \max_{j=1,...,n} \left| \frac{1}{\alpha_j} \xi_j \right|$ for any $\xi \in \mathbb{C}^n$.

We can prove easily that an element of $\mathcal{V}_\alpha^{\mathbb{C}(x)}$ defines an analytic function on $D_{a,C} := \bigcup_{K > 0, \varepsilon > 0} \bigcup_{\varepsilon < K^a} C_{K,\varepsilon}$

where

$C_{K,\varepsilon} := \{ x \in \mathbb{C}^n / d_\alpha(x, \theta^{-1}(0)) > K ||x||_\alpha \text{ and } ||x||_\alpha < \varepsilon \}$.

Here is a picture showing an example of such domain (in grey) for $n = 2$ and $\alpha = (1, 1)$:

This remark is the key point for proving the following theorem:

**Theorem 4.3.14.** — [To90][Ron13b] Let $k$ be a characteristic zero field and $\alpha \in \mathbb{R}^n_{>0}$. Let $P(Z) \in k[[x]][Z]$ be a monic polynomial whose discriminant is equal to $\delta u$ where $\delta \in k[x]$ is $(\alpha)$-homogeneous and $u \in k[[x]]$ is a unit. If $P(Z) = P_1(Z)...P_s(Z)$ where the $P_i(Z)$’s are irreducible monic polynomials of $k[[x]][Z]$, then the $P_i(Z)$’s remain irreducible in $\mathcal{V}_\alpha[Z]$.

**Sketch of proof.** — We prove the theorem in several steps:

- **First step:** Let us assume that $k = \mathbb{C}$, $\alpha \in \mathbb{N}^n$ and the coefficients of $P(Z)$ are convergent power series. Let $Q(Z)$ be an irreducible monic factor of $P(Z)$ in $\mathcal{V}_\alpha[Z]$. By Theorem 4.3.13 and the remark that follows this theorem the coefficients of $Q(Z)$ are in $\mathcal{V}_\alpha^{\mathbb{C}(x)}$. Thus the coefficients of $Q(Z)$ define analytic functions on a domain $D_{a,C}$. We can shrink $D_{a,C}$ in order to assume that
\[\delta^{-1}(0) \cap D_{a,C} = \emptyset.\] On the other hand the coefficients of \(Q(Z)\) are polynomials depending on the roots of \(P(Z)\) which are locally analytic functions outside of \(\delta^{-1}(0)\). Thus the coefficients of \(Q(Z)\) are analytic on \(D_{a,C}\) and locally analytic outside of \(\delta^{-1}(0)\). Since any point outside of \(\delta^{-1}(0)\) can be moved to a point of \(D_{a,c}\) along a curve parametrized as follows: \(t \mapsto (c_1 t^{\alpha_1}, \ldots, c_n t^{\alpha_n})\), the monodromy Theorem asserts that the coefficients of \(Q(Z)\) are in fact analytic outside \(\delta^{-1}(0)\). Since \(P(Z)\) is monic, its roots are bounded near the origin, hence the coefficients of \(Q(Z)\) also. Thus these may be extended to analytic functions in a neighborhood of the origin. Thus the coefficients of \(Q(Z)\) are analytic and \(Q(Z) = P_i(Z)\) for some \(i\) since the \(P_i(Z)\)'s are irreducible in \(\mathbb{C}\{x\}\).

• Second step: Now we prove the case \(k = \mathbb{C}, \alpha \in \mathbb{R}_{>0}\) and the coefficients of \(P(Z)\) are in \(\mathbb{C}\{x\}\). This can be done by approximating \(\nu_\alpha\) by monomial valuations whose weights are positive integers. This part is a bit technical, so we do not give more details here, but this can be done thanks to the particular form of the elements of \(V_\alpha\).

• Third step: The general case is proven in a similar way as the general case of the proof of Theorem 4.3.3 (see last step of this proof): we embed \(k\) in \(\mathbb{C}\) and we use Artin approximation Theorem. Then we conclude by using the following proposition that allows us to come back to \(k[x]\) (we skip the details here):

**Proposition 4.3.15.** — [Ron13b] Let \(k \rightarrow k'\) be a field extension. Let \(f \in k'[x]\) be a power series which is algebraic over \(k[x]\) and let \(L\) be the extension of \(k\) generated by the coefficients of \(f\). Then \(k \rightarrow L\) is a finite field extension.

This proposition is proven by using Theorem 4.3.13 and generalizes the main theorem of [CuKa08] in characteristic zero.

**Remark 4.3.16.** — In fact we can show that if \(P(Z) \in k[[x]]\) is an irreducible monic polynomial satisfying the hypothesis of the previous theorem then its Galois group is isomorphic to the Galois group of the minimal polynomial of one integral homogeneous element with respect to \(\nu\). In the case of one monomial valuation whose weights are rational numbers this means that the Galois group of \(P(Z)\) is isomorphic to the Galois group of one weighted-homogeneous polynomial (cf. Remark 7.6 [Ron13b]).

We can push the previous proof a bit further in order to obtain the following result:

**Theorem 4.3.17.** — [Ron13b] Let \(k\) be a characteristic zero field and \(\alpha \in \mathbb{R}_{>0}\). Let \(P(Z) \in k[[x]]\) be a monic polynomial whose discriminant is equal to \(\delta u\) where \(\delta \in k[x]\) is \((\alpha)\)-homogeneous and \(u \in k[[x]]\) is a unit. We set
N := \text{dim}_Q(Q\alpha_1 + \cdots + Q\alpha_n). Then there exist integral homogeneous elements with respect to \nu, \gamma_1, \ldots, \gamma_N, and a (\alpha_i)-homogeneous polynomial c(x) \in k[x] such that the roots of P(Z) are in \frac{1}{c(x)}k'[x][\gamma_1, \ldots, \gamma_N] where k \rightarrow k' is a finite field extension.

This result is a generalization of the Abhyankar-Jung Theorem. Indeed, if the \alpha_i’s are \mathbb{Q}-linearly independent then the only integral homogeneous elements are the fractional powers of the x_i’s and their products. The Abhyankar-Jung Theorem corresponds exactly to the statement of Theorem 4.3.17 in this case (with the fact that c(x) may be chosen equal to one, which is quite easy to prove in this case).

Let us finish this section by mentioning the following diophantine result that gives a necessary condition for an element of \hat{k}_\nu to be algebraic over k((x)). This is an easy corollary of Corollary 3.3.31 stated in the first part of this thesis.

**Theorem 4.3.18.** — [Ron13b] Let \nu be an Abhyankar valuation and let z \in \hat{k}_\nu be algebraic over k((x)). Then there exist two constants C > 0 and \alpha \geq 1 such that

\[ |z - \frac{f}{g}|_{\nu} \geq C |g|^\alpha_{\nu} \quad \forall f, g \in F_n. \]
• (with H. Hauser) Let us consider the following equation where \( h(x) \) and the \( f_i(x) \in \mathbb{C}\{x_1, \ldots, x_n\} \) are convergent power series,

\[
\sum_{i=1}^{p} f_i(x) y_i = h(x)
\]

and let \( \hat{y}_i(x) \in \mathbb{C}[x_j, j \in J_i] \), \( 1 \leq i \leq p \), be a solution of this equation where the \( J_i \)'s are subsets of \( \{1, \ldots, n\} \) (i.e. \( \hat{y}_i(x) \) depends only on \( x_j \) with \( j \in J_i \)).

The problem is to give conditions to insure that this kind of equations has convergent power series solutions \( \tilde{y}_i(x) \) satisfying \( \tilde{y}_i(x) \in \mathbb{C}\{x_j, j \in J_i\} \) for all \( i \) (this is a kind of Artin approximation property with constraints). The idea is to show that this approximation property is satisfied if the diagram of initial exponents of the \( \mathbb{C} \)-vector space \( \sum_{i=1}^{p} \mathbb{C}[x_j, j \in J_i].f_i(x) \) is finitely generated.

One application would be the following statement:

**Statement.** — Let \( k \) be a valued field. We set \( x' := (x_1, \ldots, x_{n-1}) \) and \( x = (x_1, \ldots, x_n) \). Let \( A := k[x'] \) and \( B := k[x, x_n] / (x_1 - x_2 x_n) \). Let \( f \in k\{x\} \) with \( f(0) = 0 \) be an element whose image in \( B \) is integral over \( A \) of degree \( d \). Then \( f \) is integral over \( k\{x'\} \).

This statement is equivalent to Gabrielov Theorem (when char(\( k \)) = 0). Our goal is to prove this result over any valued field \( k \) of any characteristic. In this case we would have to use the finiteness of the sequence of key polynomials associated to an extension of an Abhyankar valuation [Te13].

Another problem is to understand when a \( \mathbb{C} \)-vector space \( E \) defined as above, \( E = \sum_{i=1}^{p} \mathbb{C}[x_j, j \in J_i].f_i(x) \), has a finitely generated diagram of initial exponents. In particular what happens when the \( f_i \)'s are polynomials or algebraic power series?
• (with H. Hauser) One problem is to describe the set of power series solutions of a system of analytic equations. The Gringberg-Kazhdan-Drinfeld Theorem asserts that the formal neighborhood of an arc on an algebraic variety (i.e. a one variable power series solution of the equations defining the variety) is isomorphic to the direct product of an (finite dimensional) algebraic variety and a smooth infinite dimensional variety (cf. \textcite{GrKa00} and \textcite{Dr02}). We would like to give a similar description in the case of power series in several variables.

We saw in Example \textcite{3.1.10} that the Galligo-Grauert-Hironaka division Theorem applies to formal power series or convergent power series. For algebraic power series, this division theorem is no more valid as shown by the Kashiwara-Gabber Example (see Example \textcite{3.5.4}). Hironaka raised the problem of characterizing the smallest class of power series stable by division. Let us consider the following statement:

**Statement.** — Let \( f, g_1, \ldots, g_s \in \mathbb{k}(x) \) be algebraic power series over a field \( \mathbb{k} \). There exists a constant \( C > 0 \) such that the following holds: Let \( r \) be the remainder of the division of \( f \) by \( g_1, \ldots, g_s \). Let us write \( r = \sum_{k=1}^{\infty} r_{n(k)} \) where \( r_h \) is a non-zero homogeneous polynomial of degree \( h \) and \((n(k))_k\) is an increasing sequence of integers. Then

\[
    n(k + 1) \leq C n(k) \quad \forall k.
\]

We are able to prove this statement in the case the ideal generated by the \( g_i \)'s is a radical ideal or a principal ideal using a very nice result of Izumi \textcite{Iz98}. Our goal is to prove this result without assumption on the ideal \((g_1, \ldots, g_s)\). In particular this statement asserts that the example of Kashiwara-Gabber is the worst example that we may obtain by dividing algebraic power series by algebraic power series.

• Theorem \textcite{3.3.31} and the fact that Theorem \textcite{3.3.1} corresponds to a Łojasiewicz inequality as in Remark \textcite{3.3.6} are very nice examples of results valid in Diophantine geometry over algebraic number fields. The problem in the Strong Artin approximation Theorem is that we restrict to the \( m_A \)-adic valuation of the local ring \( A \). But there is no canonical valuation on \( A := \mathbb{k}[t_1, \ldots, t_m] \) when \( n \geq 2 \) and in the Strong Artin approximation Theorem it would be natural to take into account others valuation than the \( m_A \)-adic valuation. Exactly as in Diophantine geometry where all the places are taken into account through heights, we would like to find a way of defining an analogue of the Artin function which depends on all (Abhyankar?) valuations centered at the maximal ideal of \( A \). Let us consider the following statement:
**Statement.** — Let \( f \in k[t_1, \ldots, t_m, X_1, \ldots, X_n] \). Then there exist two constants \( K > 0 \) and \( a > 0 \) such that

\[
\sup_{\nu} \left\{ \frac{|f(\zeta)|_{\nu}}{d_{\nu}(\zeta, f^{-1}(0))^a} \right\} \geq K, \quad \forall \zeta \in A^m
\]

where \( \nu \) runs over all divisorial valuations centered at the maximal ideal of \( A \).

We are able to prove this statement when \( k \) is an algebraically closed of characteristic zero and \( m = 2 \) in the following two cases: either \( f \in k[X_1, \ldots, X_n] \) is a binomial, either \( f \in k[X_1, X_2] \) is irreducible (we can mention that the tools used for proving this statement in both cases uses resolution of singularities in a very similar way to [Ja00]. In particular it shows that, by taking into account all the divisorial valuations, we get a kind of Lojasiewicz inequality which is not valid if we consider only the \( m_A \)-adic valuation. This statement is not very satisfactory for several reasons but we would like to find a way of taking into account all the divisorial valuations of \( A \) for defining a notion of distance in \( A^n \) that would yield a kind of linear strong Artin approximation Theorem.

• A question that seems to have relations with the previous one is the following: how to give a valuative description of the Galois group of an irreducible polynomial with coefficients in \( k[[x_1, \ldots, x_n]] \) where \( \text{char}(k) = 0 \)? We saw how to construct an algebraically closed field \( \mathcal{E}_{\nu} \) containing \( k((x)) \) for any Abhyankar valuation \( \nu \). Let us denote by \( \hat{K}_{\nu}^{\text{alg}} \) the algebraic closure of \( k((x)) \) in \( \mathcal{E}_{\nu} \) and by \( K_{\nu}^{\text{alg}} \) the algebraic closure of \( k((x)) \) in \( \mathcal{E}_{\nu} \). Thus the field extension \( k((x)) \to \mathcal{E}_{\nu}^{\text{alg}} \) of \( k((x)) \) in its algebraic closure splits into two extensions \( k((x)) \to \hat{K}_{\nu}^{\text{alg}} \to K_{\nu}^{\text{alg}} \) (see [13]). Since \( K_{\nu}^{\text{alg}} \) is an algebraic closure of \( k((x)) \) for any \( \nu \), these fields are all isomorphic. The problem is to understand the image of the Galois group \( \text{Gal}(K_{\nu}^{\text{alg}}, K_{\mu}^{\text{alg}}) \) in the Galois group \( \text{Gal}(K_{\nu}^{\text{alg}}, k((x))) \) where \( \nu \) and \( \mu \) are two Abhyankar valuations. There are good hints for expecting that \( \text{Gal}(K_{\nu}^{\text{alg}}, k((x))) \) is generated by all the \( \text{Gal}(K_{\nu}^{\text{alg}}, \hat{K}_{\nu}^{\text{alg}}) \) when \( \nu \) runs over all divisorial valuations centered at the maximal ideal in the case \( n = 2 \). This would be a kind of Hasse principle for monic polynomials with coefficients in \( k[[x_1, x_2]] \).

• We would like to compute a sharp bound of the Artin function of a cusp \( X_1^p - X_2^q \) seen as a polynomial with coefficients in \( \mathbb{C}[[t_1, t_2]] \) (with \( p \wedge q = 1 \). The situation is more clear by Theorem 4.1.5 since we know that the difficulty occurs for approximating solutions whose order is large.
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CHAPTER 5. PERSPECTIVES


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